

Stability for the lens rigidity problem

Gang Bao*

Hai Zhang†

January 16, 2014

Abstract

Let g be a Riemannian metric for \mathbb{R}^d ($d \geq 3$) which differs from the Euclidean metric only in a smooth and strictly convex bounded domain M . The lens rigidity problem is concerned with recovering the metric g inside M from the corresponding scattering relation on the boundary ∂M . In this paper, the stability of the lens rigidity problem is investigated for metrics which are a priori close to a given non-trapping metric satisfying “strong fold-regular” condition. A metric g is called strong fold-regular if for each point $x \in M$, there exists a set of geodesics passing through x whose conormal bundle covers T_x^*M . Moreover, these geodesics contain either no conjugate points or only fold conjugate points with a non-degeneracy condition. Our main result gives the first stability result for the lens rigidity problem in the case of anisotropic metrics. The approach is based on the study of the linearized inverse problem of recovering a metric from its induced geodesic flow, which is a weighted geodesic X-ray transform problem for symmetric 2-tensor fields. A key ingredient is to show that the kernel of the X-ray transform on symmetric solenoidal 2-tensor fields is of finite dimension. It remains open whether the kernel space is trivial or not.

1 Introduction and statement of the main result

Let (M, g) be a compact Riemannian manifold with boundary ∂M . Let $\mathcal{H}^t(g)$ (or \mathcal{H}^t for the ease of notation) be the geodesic flow on the tangent bundle TM and let SM be the unit tangent bundle. Denote

$$\begin{aligned} S_+ \partial M &= \{(x, \xi) : x \in \partial M, \xi \in T_x M, |\xi|_g = 1, \langle \xi, \nu(x) \rangle > 0\}; \\ S_- \partial M &= \{(x, \xi) : x \in \partial M, \xi \in T_x M, |\xi|_g = 1, \langle \xi, \nu(x) \rangle < 0\}, \end{aligned}$$

where $\nu(x)$ is the unit outward normal at x and $\langle \cdot, \cdot \rangle$ stands for the inner product. We now define the scattering relation. Heuristically, the scattering relation encodes all the information about the geodesic flow one can obtain from the boundary. It contains not only the lengths of the unit speed geodesics connecting boundary points but also the entrance and exit directions.

*Department of Mathematics, Zhejiang University, Hangzhou 310027, China; Department of Mathematics, Michigan State University, East Lansing, MI 48824. The research was supported in part by the NSF grants DMS-0968360, DMS-1211292, the ONR grant N00014-12-1-0319, a Key Project of the Major Research Plan of NSFC (No. 91130004), and a special research grant from Zhejiang University. Email: drbaogang@gmail.com

†Department of Mathematics and its Application, Ecole Normale Supérieure, 45 Rue d’Ulm, Paris 75005, France. Email: zh.hai84@gmail.com

We first introduce it in the simplest form. Assume that the metric g is known for all points on the boundary. For each $(x, \xi) \in S_- \partial M$, define $L(g)(x, \xi) > 0$ to be the first positive moment at which the unit speed geodesic passing through (x, ξ) hits the boundary ∂M . If $L(g)(x, \xi)$ does not exist, we define formally $L(g)(x, \xi) = \infty$ and call the corresponding geodesic trapped. We define $\Sigma(g) : S_- \partial M \rightarrow \overline{S_+ \partial M}$ by

$$\Sigma(g)(x, \xi) = \mathcal{H}^{L(g)}(x, \xi). \quad (1.1)$$

The manifold (M, g) is said to be non-trapping if there exists $T > 0$ such that $L(g)(x, \xi) \leq T$ for all $(x, \xi) \in S_- \partial M$. We call the pair $(\Sigma(g), L(g))$ the scattering relation induced by the metric g .

We next generalize the above definition of the scattering relation to the case when only the restriction of the metric g to the boundary is known. Following Stefanov and Uhlmann [34], let $\kappa_\pm : S_\pm \partial M \rightarrow B(\partial M)$ be the orthogonal projection onto the open unit ball tangent bundle $B(\partial M)$. Then we redefine the scattering relation by

$$\tilde{\Sigma} = \kappa_+ \circ \Sigma \circ \kappa_- : B(\partial M) \rightarrow \overline{B(\partial M)}, \quad \tilde{L} = L \circ \kappa_-^{-1} : B(\partial M) \rightarrow [0, \infty].$$

It is easy to check that the above redefined scattering relation is invariant under diffeomorphisms of M which leave the boundary ∂M fixed. More precisely, if $\phi : M \rightarrow M$ is a diffeomorphism of M with $\phi|_{\partial M} = Id$, then the pull-back of the metric g under ϕ , denoted by ϕ^*g , has the same scattering relation as g , i.e. $(\tilde{\Sigma}(\phi^*g), \tilde{L}(\phi^*g)) = (\tilde{\Sigma}(g), \tilde{L}(g))$.

For the most general definition that removes the requirement of knowing g on ∂M , we refer to [34]. Throughout the paper, we shall restrict ourselves to the case where the metrics are known on the boundary. Hence the definition (1.1) above suffices for our purpose.

The lens rigidity problem can be formulated as recovering the metric g from its induced scattering relation. It is closely related to the inverse scattering problem for metric perturbations of the Laplacian [18, 15] in the Euclidean space, and the inverse problem of determining a metric from its induced boundary hyperbolic Dirichlet-to-Neumann map [37, 2]. In addition, it is also considered in [26] in the study of the AdS/CFT duality and holography. It is known that there is no uniqueness to the lens rigidity problem in general. The first non-uniqueness example comes from the diffeomorphisms which leave the boundary ∂M fixed. In addition, trapping of geodesics also prevents the uniqueness of the problem. We refer the readers to [13] for some counterexamples. Therefore, a natural formulation for the lens rigidity problem is as follows [37]:

Conjecture 1. *Given a compact non-trapping Riemannian manifold (M, g) with boundary ∂M , the metric g can be uniquely determined by its induced scattering relation, up to the actions of diffeomorphisms which leave the boundary fixed.*

For the lens rigidity problem, only a few uniqueness results are available. Croke [9] showed that the finite quotient space of a lens rigid manifold is lens rigid. Stefanov and Uhlmann [34] proved uniqueness up to diffeomorphisms fixing the boundary for metrics sufficiently close to a generic regular metrics. Vargo [38] showed that a class of analytic metrics are lens rigid. On the other hand, some interesting special cases of the lens rigidity problem with trapped geodesics are studied in [10, 12]. More recently, Stefanov, Uhlmann and Vasy [36] have proved the lens rigidity uniqueness and stability in the conformal class where the manifold can be foliated by strictly convex hypersurfaces. Their approach is based on the study of the local boundary

rigidity problem of determining the conformal factor of a Riemannian metric near a strictly convex boundary point from the distance function measured at pairs of points nearby. We also remark that a variant of the lens rigidity problem when the metrics are restricted to a conformal class has been investigated in [39] with a Lipschitz stability result.

To our best knowledge, no stability result on the lens rigidity problem is available for general anisotropic metrics. In this paper, we aim to address this issue and give the first result. We remark that a closely related problem for the lens rigidity problem is the boundary rigidity problem, which concerns with the unique determination of a simple metric from its induced boundary distance function, i.e. the lengths of geodesics connecting boundary points. Recall that a compact Riemannian manifold with boundary is called simple if its boundary is strictly convex and the geodesics in it have no conjugate points. It was observed by Michel [21] that the lens rigidity and boundary rigidity problem are equivalent for simple metrics. We refer to [22, 23, 24, 3, 5, 17, 4, 25, 7, 8, 11, 28, 15, 30, 20, 27, 32, 6] and the references therein for various uniqueness and stability results for the boundary rigidity problem.

Although our investigation on the stability of the lens rigidity problem is motivated and inspired by the results of Stefanov and Uhlmann [31, 32, 34] on the boundary rigidity problem and the uniqueness of the lens rigidity problem, the approach and scope are different. First, the stability result in [32] and the uniqueness results in [31, 34] are based on linearizing the operator which maps metrics to the lengths of geodesics connection boundary points, while our approach is based on linearizing the operator which maps metrics to their associated geodesic flows. There are two main advantages by using the geodesic flow: one is that it allows us to handle geodesics which have conjugate points; the other is that the linearization is more straightforward, though it requires a larger manifold to work with. We remark that the information used from the geodesic flow is in fact equivalent to the scattering relation (see Appendix A). Second, the geodesic X-ray transforms of the linearized operators in [31, 32, 34] assume no conjugate points, while our transform allows the presence of conjugate points. In addition, our transform has a variable weight as opposed to the constant weight in their transforms. Note that fold type conjugate points occur commonly on geodesics in a general Riemannian manifold, while other types of conjugate points occur less commonly. In fact, it is shown that in the set of all geodesics passing through a given point $x \in M$, the set of geodesics with fold type conjugate points has the same dimension as the whole set, and the set of geodesics with other types of conjugate points has a lower dimension. Therefore, it is necessary to consider geodesics with fold type conjugate points in the study of the X-ray transform in a general Riemannian manifold [2]. In order to control the effects of fold conjugate points, we introduce the concept “strong fold-regular”. By allowing the usage of geodesics with strong fold-regular conjugate points in the inversion for the X-ray transform, we are able to study the X-ray transform on a quite general class of non-simple manifolds and consequently the lens rigidity problem as well.

Throughout the paper, our discussions are presented in the space \mathbb{R}^d for the ease of exposition. Let M be a strictly convex domain in \mathbb{R}^d . We required that the metrics equipped to M are identical to the Euclidean one, denoted by e , in a neighborhood of ∂M . This requirement saves us from technicalities caused from the boundary, while not losing much generality. In fact, for any Riemannian metric g to M , there exists another strictly convex domain $M_e \ni M$ with a metric g_e such that $g_e = e$ in some neighborhood of ∂M_e and $g_e|_M = g|_M$. The same is also true for metrics that are sufficiently small perturbations of g .

Our main result on the lens rigidity problem reads as follows (the proof will be given in Section 6).

Theorem 1.1. *Let g be a smooth Riemannian metric and M be a smooth and strictly convex bounded domain in \mathbb{R}^d with $d \geq 3$. Let $\hat{M} \Subset M$. Assume that: (1). the support of $g - e$ is contained in \hat{M} ; (2). the geodesic flow induced by g is non-trapping in M over time T ; (3). the metric g is strong fold-regular (see Section 3 for its definition). Then there exist a positive integer k , a finite dimensional space $\mathfrak{L} \in L^2(S(\tau_2 M))$ where $S(\tau_2 M)$ denotes the set of symmetric 2-tensor fields on M , and a finite number of smooth functions $\alpha_j \in C^\infty(S_{-\partial M})$, $j = 1, 2, \dots, N$, such that for any Riemannian metric \tilde{g} close enough to g in $C^{k, \frac{1}{2}}(M)$ and differs from e only in \hat{M} , there exists a diffeomorphism φ of M with the property that $\varphi^* \tilde{g} = e$ in a neighborhood of ∂M , and moreover*

$$\|\varphi^* \tilde{g} - g\|_{C^0(M)} \lesssim \sum_{j=1}^N \|\alpha_j(\mathcal{H}^T(\tilde{g}) - \mathcal{H}^T(g))\|_{H^{3+\lfloor \frac{d}{2} \rfloor}(S_{-\partial M})}, \quad (1.2)$$

provided $\varphi^* \tilde{g} - g \perp \mathfrak{L}$ in $L^2(S(\tau_2 M))$.

Remark 1.1. *It is shown in the proof of Theorem 1.1 that the integer k can be chosen to be $13 + 2\lfloor \frac{d}{2} \rfloor$, which may not be the sharpest one.*

Remark 1.2. *In Theorem 1.1, only partial information from the geodesic flow (or the scattering relation) is used. As a consequence, the assumptions on the background metric g can be significantly weakened in the following two ways:*

- (1). *The non-trapping condition can be replaced by the condition that only the geodesics which are used for the stability estimates (or those which are in the support of the cut-off functions α_j 's) are non-trapping in M for time T .*
- (2). *The convexity condition on the boundary ∂M can be replaced by the condition that the geodesics used are transverse to the boundary ∂M .*

Remark 1.3. *In Theorem 1.1, if we have another metric \tilde{g}_1 which satisfies the same conditions as \tilde{g} , then the conclusion becomes*

$$\|\varphi^* \tilde{g} - \varphi_1^* \tilde{g}_1\|_{C^0(M)} \lesssim \sum_{j=1}^N \|\alpha_j(\mathcal{H}^T(\tilde{g}) - \mathcal{H}^T(\tilde{g}_1))\|_{H^{3+\lfloor \frac{d}{2} \rfloor}(S_{-\partial M})},$$

provided $\varphi^* \tilde{g} - \varphi_1^* \tilde{g}_1 \perp \mathfrak{L}$, where φ_1 is determined by \tilde{g}_1 in the same way as φ by \tilde{g} .

Remark 1.4. *The finite dimensional space \mathfrak{L} is related to the kernel of the linearized inverse problem of determining g from its induced geodesic flow $\mathcal{H}^T(g)|_{S_{-\partial M}}$, or more precisely, the geodesic X-ray transform \mathfrak{X} (see Section 2.5) for solenoidal 2-tensor fields. It is not clear that whether \mathfrak{L} is trivial or not for a general non-simple metric. This problem is usually referred to as the “s-injectivity problem” in the literature, see for instance [31, 33, 16].*

Remark 1.5. *Theorem 1.1 can be stated in terms of the scattering relation instead of the geodesic flow. In fact, there are two ways to do so. One way is to replace the geodesic flow by the corresponding scattering relation with the help of Identity (6.10). The other way is to linearize the scattering relation with respect to the metric as for the geodesic flow. We can show that a similar conclusion holds if we replace \mathcal{H}^T by (Σ, L) .*

We now give a brief account of the main approach in the paper. We first linearize the operator which associates a metric with its induced geodesic flow on the cosphere bundle at a fixed smooth background metric g . We assume that g is non-trapping in M over the time T . This leads to an X-ray transform operator \mathfrak{X} for symmetric 2-tensor fields. We form the normal operator $\mathfrak{N} = \mathfrak{X}^\dagger \mathfrak{X}$. In the case when the metric g is not simple (or the exponential maps of the manifold (M, g) have singularities), the Schwartz kernel of \mathfrak{N} has two types of singularities: one is from the diagonal, the other is from conjugate points. We show that the former singularities yield a Ψ DO (pseudo differential operator) while the latter ones yield a FIO (Fourier integral operator) locally in the case when the singularities are of fold type.

We study the stability of the above X-ray transform operator \mathfrak{X} for symmetric 2-tensor fields. Since it vanishes for potential 2-tensor fields, we restrict to the solenoidal 2-tensor fields. By micro-local analysis, a set of geodesics passing through x whose conormal bundle can cover the cotangent space $T_x^*\mathbb{R}^d$ are needed in order to reconstruct the solenoidal part of a 2-tensor field f at x from its X-ray transform $\mathfrak{X}f$. We allow fold conjugate points along these geodesics, but require that these conjugate points contribute to a smoother term than the point x itself. This is the case when the intrinsic derivatives of the differential of the exponential map at the conjugate vectors satisfy a non-degenerate condition, which is called “strong fold-regular”. We establish a local stability estimate near “strong fold-regular” points. These local estimates are glued together to form a global one. We remark that it remains open whether the X-ray transform operator has a non-trivial kernel in the space of solenoidal 2-tensor fields. This problem is often referred to as the “s-injectivity” problem. By imposing some orthogonality conditions, we obtain a Lipschitz type stability estimate (see Theorem 3.2).

We remark that a similar X-ray transform operator is obtained by linearizing the operator which maps metrics within a conformal class to their induced geodesic flows at a smooth background metric in [2]. The transform there is applied to scalar functions. Its kernel can be shown to be of finite dimension. In comparison, the X-ray transform \mathfrak{X} in this paper has an infinitely dimensional kernel space. In fact, it vanishes for all potential 2-tensor fields as shown in Lemma 5.1. Thus the analysis for its invertibility and stability is much more subtle here.

The Lipschitz type estimate for the X-ray transform \mathfrak{X} , or the linearized inverse problem of recovering a metric from its induced geodesic flow, enables us to derive a Lipschitz stability result for the nonlinear inverse problem. This last step is made possible by using the construction for diffeomorphisms in [34].

The paper is organized as follows: In Section 2, we present some preliminaries and introduce two X-ray transform operators \mathfrak{I} and \mathfrak{X} which are obtained from linearizing the operator which maps a metric to its induced geodesic flow. The results on the stability estimates for the X-ray transform \mathfrak{X} are stated in Section 3. In Section 4, we derive properties of the X-ray transform \mathfrak{I} and its normal \mathfrak{M} . Based on these properties, we study the X-ray transform \mathfrak{X} in Section 5 and prove the stability results stated in Section 3. Finally, we prove the main result of the paper, Theorem 1.1, in Section 6.

2 Preliminaries

2.1 Notations

Throughout the paper, we denote by $\exp_x(\cdot) : T_x\mathbb{R}^d \rightarrow \mathbb{R}^d$ the exponential map at $x \in \mathbb{R}^d$ with respect to the background metric g . The differential of $\exp_x(\cdot)$ at $\xi_* \in T_x\mathbb{R}^d$ is denoted by

$d_\xi \exp_x(\xi_*)$. The norm of $\xi \in T_x \mathbb{R}^d$ is denoted by $|\xi|$. If the linear map $d_\xi \exp_x(\xi_*)$ is not a diffeomorphism, we call the vector ξ_* a conjugate vector, and the corresponding point $\exp_x(\xi_*)$ a conjugate point. The kernel of the linear map $d_\xi \exp_x(\xi_*)$ is denoted by $N_x(\xi_*)$, which is a subspace of $T_x \mathbb{R}^d$.

In addition, we introduce the following conventions:

1. Let A_1 and A_2 be two matrices (including vectors which can be regarded as single column or single row matrices), then the product of A_1 and A_2 is denoted by $A_1 \cdot A_2$. Sometimes, the dot is omitted for simplicity;
2. Let A be a complex matrix, then A^\dagger stands for its conjugate transpose. If A is a linear operator in a Hilbert space, then A^\dagger stands for its formal adjoint. If A is real and symmetric and C a real number, then $A \geq C$ means that the matrix $A - C \cdot Id$ is symmetric and positive definite.
3. Let U and V be two open set in a metric space, then $U \Subset V$ means that the closure of U , denoted by \bar{U} is compact and is a subset of V ;
4. Let C_1 and C_2 be two positive numbers, then $C_1 \lesssim C_2$ means that $C_1 \leq C \cdot C_2$ for some constant $C > 0$ independent of C_1 and C_2 .

2.2 Symmetric tensor fields

We denote by $S(\tau_2 M)$ the set of symmetric covariant 2-tensor fields to the manifold M . In the natural coordinate of \mathbb{R}^d , each $f \in S(\tau_2 M)$ is assigned a family of functions $\{f_{ij}; 1 \leq i, j \leq d\}$ in M such that $f_{ij} = f_{ji}$. We denote by $S(\tau_2 M, \mathbb{R}^d)$ the set of families of functions $\{\Pi_{ij}^k; 1 \leq i, j, k \leq d\}$ in M such that $\Pi_{ij}^k = \Pi_{ji}^k$. An example of $S(\tau_2 M, \mathbb{R}^d)$ is the set of Christoffel symbols. We remark that elements in $S(\tau_2 M, \mathbb{R}^d)$ are not tensor fields, i.e. they are not invariant under the change of coordinates. We also remark that we shall use the natural coordinate of \mathbb{R}^d throughout the paper.

Let $L^2(S(\tau_2 M))$ be the Hilbert space of L^2 -integrable symmetric covariant 2-tensor fields with the following inner product

$$(f, h)_{L^2(M)} = \int_M f_{ij} h^{ij} \sqrt{\det g} dx = \int_M f_{ij} g^{ii'} g^{jj'} h_{i'j'} \sqrt{\det g} dx, \quad (2.1)$$

where $\{h^{ij}\}$ is the contravariant tensor field corresponding to $\{h_{ij}\}$.

Similar to the above L^2 space for 2-tensor fields, we can define L^2 space for any tensor fields. Here we remark that throughout the paper, the Einstein's summation rule is applied whenever there are sub and superscripts with the same label. For convenience, we always use the usual convection of raising and lowering indices and we treat the covariant and the corresponding contravariant tensor field as two representations of the same tensor field.

Based on the L^2 space defined above, one can define the H^k space for any positive integer k by the following inductive formula:

$$(f, h)_{H^k(M)} = (f, h)_{L^2(M)} + (\nabla f, \nabla h)_{H^{k-1}(M)}.$$

By the standard interpolation theory, one can define the H^s space for any $s \geq 0$ as well.

Remark 2.1. One can define an equivalent norm for the space of 2-tensor fields $L^2(S(\tau_2 M))$ by using the natural coordinate of \mathbb{R}^d . Indeed, for $f, h \in L^2(S(\tau_2 M))$, define

$$(f, h)_{\tilde{L}^2(M)} = \sum_{ij} \int_M f_{ij}(x) h_{ij}(x) dx. \quad (2.2)$$

We can check the equivalence of the two norms by using the open mapping theorem. Similarly, we can also define the equivalent norm for the space $H^k(S(\tau_2 M))$.

The norm induce by the inner product (2.2) shall be used frequently in the subsequent analysis, since it is more convenient to work with than the one induced by (2.1). Considering that the two norms are equivalent, we use the same symbol $\|\cdot\|_{L^2}$ (or $\|\cdot\|_{H^k}$) to denote them for simplicity.

We define $L^2(S(\tau_2 M), \mathbb{R}^d)$ to be the Hilbert space of L^2 -integrable fields in $S(\tau_2 M, \mathbb{R}^d)$ with the following inner product

$$(\Pi, \Gamma) = \int_M \Pi_{ij}^k \Gamma_k^{ij} \sqrt{\det g} dx,$$

where $\Gamma_k^{ij} = \delta_{kk'} g^{ii'} g^{jj'} \Gamma_{i'j'}^{k'}$ with $\delta_{kk'}$ being the standard Cronecker's symbol. We define the H^k space for $S(\tau_2 M, \mathbb{R}^d)$ by using the inner product (2.2) for $k \geq 1$.

Let \tilde{M} be a compact set in M and let k be a nonnegative integer, we denote by $H^k(S(\tau_2(\tilde{M}, M)))$ the subspace of $H^k(S(\tau_2 M))$ whose elements are supported in \tilde{M} .

Finally, we introduce a lifting operator. For each $h \in S(\tau_2 M, \mathbb{R}^d)$, we define $\iota(h) \in S(\tau_2 M, \mathbb{R}^{2d})$ by setting the first d components to be zero.

2.3 Decomposition of symmetric 2-tensor fields

We present some basic facts about the decomposition of symmetric 2-tensor fields. We refer to [29] for details. Given a symmetric 2-tensor field $f = \{f_{ij}\}$, we define the 1-tensor field $\delta^s f$ by

$$[\delta^s f]_i = g^{jk} \nabla_k f_{ij},$$

where ∇_k is the covariant derivative. On the other hand, given a 1-tensor field $v = \{v_i\}$, we define the 2-tensor field $d^s v$, called symmetric differential of v , by

$$[d^s v]_{ij} = \frac{1}{2} (\nabla_i v_j + \nabla_j v_i).$$

The operators d^s and $-\delta^s$ are formally adjoint to each other in the L^2 space.

For each $f \in L^2(S(\tau_2 M))$, there exists a unique orthogonal decomposition

$$f = f^s + d^s v = \mathcal{S}f + \mathcal{P}f,$$

such that $f^s \in L^2(S(\tau_2 M))$ satisfies $\delta^s f^s = 0$ and $v \in H_0^1(M)$ (i.e. $v = 0$ on ∂M). The fields f^s and $d^s v$ are called the solenoidal and potential parts of f , respectively. The following orthogonality property holds for the two projectors $\mathcal{S}, \mathcal{P} : L^2(S(\tau_2 M)) \rightarrow L^2(S(\tau_2 M))$:

$$\mathcal{P}\mathcal{S} = \mathcal{S}\mathcal{P} = 0.$$

The decomposition can be constructed as follows. Consider the operator $\Delta_M^s = \delta^s d^s$. It

can be verified that Δ_M^s is elliptic in M and the corresponding Dirichlet problem satisfies the Lopatinskii condition. Thus Δ_M^s is invertible with the Dirichlet boundary condition. Denote by $(\Delta_M^s)^{-1}$ its inverse. Then

$$v = (\Delta_M^s)^{-1} \delta^s f, \quad f^s = f - d^s (\Delta_M^s)^{-1} \delta^s f.$$

2.4 Decomposition of symmetric 2-tensors

We introduce a similar decomposition for symmetric tensors (see Section 2.6 in [28] for detail). Denote

$$S_2 = \{\{v_{ij}\}; v_{ij} = v_{ji} \in \mathbb{R}, 1 \leq i, j \leq d\}; \quad S_1 = \mathbb{R}^d.$$

For each $x \in \mathbb{R}^d \setminus 0$, we define $i_x : S_1 \rightarrow S_2$ and $j_x : S_2 \rightarrow S_1$ in the following way:

$$(i_x v)_{ij} = \frac{1}{2} (v_i x_j + v_j x_i), \quad (j_x u)_i = u_{ij} x^j,$$

where $v \in S_1$ and $u \in S_2$. The two operators i_x and j_x are called the symmetric multiplication and convolution with x , respectively. We can show that for every $f \in S_2$ and $x \in \mathbb{R}^d \setminus 0$, there exist uniquely determined $h \in S_2$ and $v \in S_1$ such that

$$f = h + i_x v, \quad j_x h = 0.$$

2.5 Linearization of the geodesic flow with respect to the metric

We linearize the operator which maps a metric to its induced geodesic flow restricted to the cosphere bundle. Recall that g is the fixed background metric. Denote by $\Gamma = \{\Gamma_{ij}^k; 1 \leq i, j, k \leq d\}$ its Christoffel symbol. Let $f \in C^\infty(S(\tau_2 M))$ be such that its support is contained in $\hat{M} \Subset M$. Then $g + f$ defines a perturbed metric in M provided $\|f\|_{C^2(M)} \ll 1$. Denote by $\mathcal{H}^T(g + f)$ the corresponding geodesic flow. We now derive an explicit formula for the linearized map $\frac{\delta \mathcal{H}^T(g)}{\delta g}$ at g below.

We first define a matrix $\Phi(x, \xi)$ for each (x, ξ) that lies on the orbit of the set $S_- \partial M$ under the geodesic flow $\mathcal{H}^t(g)$ ($t \geq 0$), see also [2]. Let $(x_0, \xi_0) = \mathcal{H}^{L^-(g)(x, \xi)}(g)(x, \xi) \in S_- \partial M$, where $L^-(g)(x, \xi)$ is the first negative moment the geodesic orbit $\mathcal{H}^t(g)(x, \xi)$ hits the boundary $S_- \partial M$. Let $\phi(t)$ be the solution of the following ODE system

$$\dot{\phi}(t) = -\phi(t) A(\mathcal{H}^t(g)(x_0, \xi_0)), \quad \phi(0) = Id,$$

where $A(x, \xi) = (\frac{\partial H}{\partial x}, \frac{\partial H}{\partial \xi})$ is a $2d \times 2d$ matrix and H is given by

$$H(x, \xi) = (\xi^1, \xi^2, \dots, \xi^d, -\Gamma_{ij}^1 \xi^i \xi^j, -\Gamma_{ij}^2 \xi^i \xi^j, \dots, -\Gamma_{ij}^d \xi^i \xi^j)^\dagger.$$

The matrix $\Phi(x, \xi)$ is defined by

$$\Phi(x, \xi) = \phi(T)^{-1} \phi(-L^-(g)(x_0, \xi_0)).$$

It is clear that $\Phi(\cdot, \cdot)$ is well-defined and is smooth in a sufficiently small neighborhood of any unit speed geodesic orbit if only the corresponding geodesic is transverse to the boundary ∂M .

With the matrix-valued weight Φ defined above, we now define the X-ray transform operator $\mathfrak{I} : C^\infty(S(\tau_2 M); \mathbb{R}^{2d}) \rightarrow \mathcal{D}'(S_{-\partial M}; \mathbb{R}^{2d})$ by

$$\begin{aligned} (\mathfrak{I}\Pi)(x_0, \xi_0) &= \int_0^T \Phi(x(s), \xi(s)) \Pi_{ij}(x(s)) \xi^i(s) \xi^j(s) ds \\ &= \int_0^{L(g)(x_0, \xi_0)} \Phi(x(s), \xi(s)) \Pi_{ij}(x(s)) \xi^i(t) \xi^j(s) ds \end{aligned} \quad (2.3)$$

where $(x(s), \xi(s)) = \mathcal{H}^s(g)(x_0, \xi_0)$ and $\Pi_{ij} \xi^i \xi^j$ is viewed as a vector in \mathbb{R}^{2d} with the k -th component given by $\Pi_{ij}^k \xi^i \xi^j$.

To obtain $\frac{\delta \mathcal{H}^T(g)}{\delta g}$, we need to introduce one more operator. For each $f \in S(\tau_2 \mathbb{R}^d)$, we define

$$(\mathcal{L}f)_{ij}^k = \frac{1}{2} g^{lk} \left(\frac{\partial f_{jl}}{\partial x^i} + \frac{\partial f_{il}}{\partial x^j} - \frac{\partial f_{ij}}{\partial x^l} \right) - g^{lk} \Gamma_{ij}^m f_{ml}.$$

It is clear that \mathcal{L} is a first order partial differential operator. We denote by \mathcal{L}^\dagger is formal adjoint.

Denote $\mathfrak{X} = \mathfrak{I} \circ \iota \circ \mathcal{L}$. Then we have $\frac{\delta \mathcal{H}^T(g)}{\delta g}|_{S_{-\partial M}} = \mathfrak{X}$. In fact, the following result is proved in Appendix B.

Proposition 1. *The following estimate holds*

$$\|\mathcal{H}^T(g + f) - \mathcal{H}^T(g) - \mathfrak{X}f\|_{C^k(S_{-\partial M})} \lesssim \|f\|_{C^{k+2}(M)}^2$$

for any nonnegative integer k .

Remark 2.2. *Formula (2.3) is derived in the coordinate of $T\mathbb{R}^d$. It may not be geometrically invariant.*

3 Stability for the geodesic X-ray transform \mathfrak{X}

Recently, “fold-regular” metrics have been introduced in [2] to study the geodesic X-ray transform for scalar functions in a general Riemannian manifold. They generalize the “regular” metrics in [33] by allowing fold conjugate points along the geodesics which are used in the inversion of the X-ray transform. In this section, we present some stability results on the transform \mathfrak{X} for 2-tensor fields obtained in the previous section under the “strong fold-regular” condition. We refer to [28, 29, 31, 32] and the most recent work [16] and the reference therein for the study on X-ray transforms in a simple manifold and their applications.

We first introduce some definitions. We refer to [2] for discussions on these concepts. Recall that a conjugate vector $\xi_* \in T_x \mathbb{R}^d$ is of fold type if the following two conditions are satisfied: (1) the rank of the linear mapping $d_\xi \exp_x(\xi_*)$ equals to $d - 1$ and the function $\det(d_\xi \exp_x(\xi))$ vanishes of order 1 at ξ_* ; (2) the kernel space $N_x(\xi_*)$ of the linear mapping $d_\xi \exp_x(\xi_*)$ is transversal to the manifold $\{\eta : \det(d_\eta \exp_x(\eta)) = 0\}$ at ξ_* .

Definition 3.1. *A fold vector $\xi_* \in T_x M$ is called strong fold-regular if the following condition is satisfied*

$$d_\xi^2 \exp_x(\xi_*)(N_x(\xi_*) \setminus 0 \times \cdot)|_{T_{\xi_*} S(x)} \quad \text{is of full rank} \quad (3.1)$$

where $S(x) = \{\xi \in T_x \mathbb{R}^d; \det(d_\xi \exp_x \xi) = 0\}$ and $T_{\xi_*} S(x)$ is its tangent space at ξ_* , and $d_\xi^2 \exp_x(\xi_*)(N_x(\xi_*) \setminus 0 \times \cdot)$ is the intrinsic derivative of the differential of the exponential map $d_\xi \exp_x(\cdot) : T_x M \rightarrow T_{\exp_x \xi} M$ at (x, ξ_*) .

We refer to [1] for the definition of the intrinsic derivative. The above condition (3.1) was first introduced in [35], which guaranties the graph condition ([19]) for the FIO obtained from the X-ray transform studied therein. We use it for the same purpose here, see Lemma 4.6.

Definition 3.2. A point $x \in M$ is called strong fold-regular if there exists a compact subset $\mathcal{Z}(x) \subset S_x M$ such that the following two conditions are satisfied:

1. For each $\xi \in \mathcal{Z}(x)$, there exist no singular vectors except those of strong fold-regular type along the ray $\{t\xi : t \in \mathbb{R}\}$ for the map $\exp_x(\cdot)$; moreover, the corresponding geodesic hits the boundary ∂M transversely.
2. $\forall \xi \in S_x M, \exists \theta \in \mathcal{Z}(x)$, such that $\theta \perp \xi$.

The second condition in the above definition can be viewed as “completeness” condition, see [33] for instance. It ensures the ellipticity of the Ψ DO part of the normal of properly truncated X-ray transforms, see Lemma 5.4.

We next introduce the truncated X-ray transform operator. For any $\alpha \in C_0^\infty(S_- \partial M)$, we define the truncated operator \mathfrak{X}_α by

$$\mathfrak{X}_\alpha f(x_0, \xi_0) = \alpha(x_0, \xi_0) \mathfrak{X}_\alpha f(x_0, \xi_0). \quad (3.2)$$

Denote by $\mathfrak{N}_\alpha = \mathfrak{X}_\alpha^\dagger \mathfrak{X}_\alpha$ the normal operator.

We now present a local stability estimate near strong fold-regular points and several corollaries. The proofs will be given in Section 5.2.

Theorem 3.1. Let $x_* \in \tilde{M} \Subset M$ be a strong fold-regular point and let k be a nonnegative integer. Then there exist two neighborhoods $U(x_*) \Subset \tilde{U}(x_*)$ of x_* , and a cutoff function $\alpha \in C_0^\infty(S_- \partial M)$, such that the following estimate holds uniformly for all $f \in H^k(S(\tau_2(\tilde{M}, M)))$

$$\|f^s\|_{H^k(U(x_*))} \lesssim \|\mathfrak{N}_\alpha f\|_{H^{k-1}(\tilde{U}(x_*))} + \|f\|_{H^{k+1-\frac{d}{2}}(M)}. \quad (3.3)$$

The above local result can be extended to a global one.

Definition 3.3. The background metric g is called strong fold-regular if all points in M are strong fold-regular with respect to the geodesic flow $\mathcal{H}^t(g)$.

Corollary 3.1. Assume that the background metric g is strong fold-regular. Let \tilde{M} be a strictly convex and smooth sub-domain of M . Let k be a nonnegative integer. Then there exist $\tilde{U}(x_j) \subset M$, $\alpha_j \in C_0^\infty(S_- \partial M)$, $j = 1, 2, \dots, N$, such that the following estimate holds for all $f \in H^k(S(\tau_2(\tilde{M}, M)))$:

$$\|f^s\|_{H^k(M)} \lesssim \sum_{j=1}^N \|\mathfrak{N}_{\alpha_j} f\|_{H^{k-1}(\tilde{U}(x_j))} + \|f\|_{H^{k+1-\frac{d}{2}}(M)}.$$

Corollary 3.2. *Assume that the background metric g is strong fold-regular. Let \tilde{M} and $\alpha_j \in C_0^\infty(S_{-\partial M})$, $j = 1, 2, \dots, N$, be as in corollary 3.1. Then there exists a finite dimensional space $\mathfrak{L} \in \mathcal{SL}^2(S(\tau_2(\tilde{M}, M)))$ such that the following estimate holds for all $f \in L^2(S(\tau_2(\tilde{M}, M)))$ satisfying $f^s \perp \mathfrak{L}$ in $L^2(S(\tau_2 M))$:*

$$\|f^s\|_{L^2(M)} \lesssim \sum_{j=1}^N \|\mathfrak{N}_{\alpha_j} f\|_{H^{-1}(\tilde{U}(x_j))}. \quad (3.4)$$

Corollary 3.3. *Assume that the background metric g is strong fold-regular. Let \tilde{M} , \mathfrak{L} and $\alpha_j \in C_0^\infty(S_{-\partial M})$, $j = 1, 2, \dots, N$, be as in Corollary 3.2, and let k be a nonnegative integer. Then the following estimate holds for all $f \in H^k(S(\tau_2(\tilde{M}, M)))$ satisfying $f^s \perp \mathfrak{L}$ in $L^2(S(\tau_2 M))$:*

$$\|f^s\|_{H^k(M)} \lesssim \sum_{j=1}^N \|\mathfrak{N}_{\alpha_j} f\|_{H^{k-1}(\tilde{U}(x_j))}. \quad (3.5)$$

Remark 3.1. *The strong fold-regular condition in Corollary 3.2 and 3.3 may be weakened. In fact, as in [2], we can define a fold vector $\xi \in T_x M$ to be fold-regular if the operator germ \mathfrak{M}_ξ , which characterizes contributions from an infinitesimal small neighborhood of the point $\exp_x \xi$ to the normal operator \mathfrak{M} , is compact from $H^k(\tilde{M}, M)$ to $H^{k+1}(U(x))$ for any integer k and some neighborhood $U(x)$ of x . Then the results in Corollary 3.2 and 3.3 and consequently Theorem 1.1 also hold under the fold-regular condition.*

4 The X-ray transform operator \mathfrak{I} and its normal \mathfrak{M}

We study the X-ray transform operator \mathfrak{I} and its normal \mathfrak{M} in this section. We first introduce some measures that are needed in the integrals we use later. We refer [29] for more detail. Let $x \in M$, we denote by $d\mu_x(\xi)$ and $d\mu(x, \xi)$ the measure induced by the metric g to the sphere $S_x M$ and the spherical bundle SM , respectively. For the set $S_{-\partial M}$ characterizing the set of geodesics passing through M , its canonical measure is given by $|\langle \nu(x), \xi \rangle| d\Sigma^{2d-2}(x, \xi)$, where $d\Sigma^{2d-2}(x, \xi)$ is the restriction of the measure $d\mu(x, \xi)$ to the subset $S_{-\partial M}$ and $\nu(x)$ is the outward normal to ∂M at x . In the coordinates of \mathbb{R}^d , we have

$$\begin{aligned} d\mu_x(\xi) &= (\det g)^{\frac{1}{2}} \sum_{i=1}^d (-1)^{i-1} \xi^i d\xi^1 \wedge \dots \wedge d\xi^{i-1} \wedge d\xi^{i+1} \wedge \dots \wedge d\xi^d; \\ d\mu(x, \xi) &= d\mu_x(\xi) \wedge (\det g)^{\frac{1}{2}} dx. \end{aligned}$$

We also use $d\mu_x(\xi)$ to denote the measure to the tangent space $T_x M$. In coordinates, $d\mu_x(\xi) = (\det g)^{\frac{1}{2}} d\xi$.

We now present some basic properties about the X-ray transform \mathfrak{I} .

Lemma 4.1. *The X-ray transform operator \mathfrak{I} is bounded from $H^k(S(\tau_2 M, \mathbb{R}^{2d}))$ to $H^k(S_{-\partial M}, \mathbb{R}^{2d})$ for any integer $k \geq 0$.*

Proof. The lemma can be proved in a similar way as Theorem 3.3.1 in [29].

Lemma 4.2. *The transpose $\mathfrak{I}^\dagger : H^k(S_{-\partial M}, \mathbb{R}^{2d}) \rightarrow H^k(S(\tau_2 M, \mathbb{R}^{2d}))$ of the operator \mathfrak{I} has the following representation:*

$$\mathfrak{I}^\dagger(h)_{ij}(x) = \int_{S_x M} \Phi^\dagger(x, \xi) h^\sharp(x, \xi) \xi_i \xi_j d\mu_x(\xi).$$

where h^\sharp is the unique lift of h to SM which is invariant under the geodesic flow and is equal to h on $S_{-\partial M}$.

Proof. For any $h \in C^\infty(S_{-\partial M}, \mathbb{R}^{2d})$, we have

$$\begin{aligned} (\mathfrak{I}(\Pi), h) &= \int_{S_{-\partial M}} |\langle \nu(x_0), \xi_0 \rangle| d\Sigma^{2d-2}(x_0, \xi_0) h(x_0, \xi_0) \int_0^\infty \Phi(\mathcal{H}^t(x_0, \xi_0)) \Pi_{ij}(x(t)) \xi^i(t) \xi^j(t) dt \\ &= \int_{SM} \langle h^\sharp(x, \xi), \Phi(x, \xi) \Pi_{ij}(x) \xi^i \xi^j \rangle d\mu(x, \xi) \quad (\text{by Santalo's formula, see [29]}) \\ &= \int_{SM} \langle \Phi^\dagger(x, \xi) h^\sharp(x, \xi), \Pi_{ij}(x) \xi^i \xi^j \rangle d\mu(x, \xi) \\ &= \int_{SM} \sum_{1 \leq i, j \leq d} \langle \Phi^\dagger(x, \xi) h^\sharp(x, \xi) \xi^i \xi^j, \Pi_{ij}(x) \rangle d\mu(x, \xi) \\ &= \left(\int_{S_x M} \Phi^\dagger(x, \xi) h^\sharp(x, \xi) \xi_i \xi_j d\mu_x(\xi), \Pi_{ij}(x) \right), \end{aligned}$$

whence the claim follows.

Lemma 4.3. *The normal operator $\mathfrak{M} = \mathfrak{I}^\dagger \mathfrak{I}$ has the following representation:*

$$\begin{aligned} \mathfrak{M}(\Pi)_{ij}(x) &= \int_{S_x M} d\mu_x(\xi) \int_{-\infty}^\infty \Phi^\dagger(x, \xi) \Phi(\mathcal{H}^t(x, \xi)) \xi_i \xi_j \Pi_{mn}(\exp_x t\xi) \cdot \exp_x^m t\xi \cdot \exp_x^m t\xi dt \\ &= \int_{T_x M} W_{ij}^{mn}(x, \xi) f_{mn}(\exp_x \xi) d\mu_x(\xi), \end{aligned}$$

where

$$\begin{aligned} W_{ij}^{mn}(x, \xi) &= \frac{\xi_i \xi_j}{|\xi|^{d+1}} \exp_x^m \left(\frac{\xi}{|\xi|} \right) \exp_x^n \left(\frac{\xi}{|\xi|} \right) \\ &\quad \cdot \left(\Phi^\dagger(x, \frac{\xi}{|\xi|}) \Phi(\exp_x \xi, \exp_x \frac{\xi}{|\xi|}) + \Phi^\dagger(x, -\frac{\xi}{|\xi|}) \Phi(\exp_x \xi, -\exp_x \frac{\xi}{|\xi|}) \right). \end{aligned}$$

We now show some local properties of the normal operator \mathfrak{M} .

Let $\tilde{M} \Subset M$. From now on, we fix $x_* \in \tilde{M}$. We first decompose \mathfrak{M} locally into two parts based on the separation of singularities of its Schwartz kernel. By the existence of uniformly normal neighborhood in Riemannian manifold, we can find $\epsilon_2 > 0$ and a neighborhood of x_* , say $\tilde{U}(x_*) \subset \mathbb{R}^d$, such that

$$\exp(x, \cdot)|_{|\xi| < 2\epsilon_2} \text{ is a diffeomorphism for any } x \in \tilde{U}(x_*). \quad (4.1)$$

Let $\chi_{\epsilon_2} \in C_0^\infty(\mathbb{R})$ be such that $\chi(t) = 1$ for $|t| < \epsilon_2$ and $\chi(t) = 0$ for $|t| > 2\epsilon_2$. We then

define

$$(\mathfrak{M}_1 f)_{ij}(x) = \int_{T_x M} W_{ij}^{mn}(x, \xi) f_{mn}(\exp_x \xi) \chi_{\epsilon_2}(|\xi|) d\mu_x(\xi), \quad (4.2)$$

$$(\mathfrak{M}_2 f)_{ij}(x) = \int_{T_x M} W_{ij}^{mn}(x, \xi) f_{mn}(\exp_x \xi) (1 - \chi_{\epsilon_2}(|\xi|)) d\mu_x(\xi). \quad (4.3)$$

Note that for any f supported in \tilde{M} , $f(\exp_x \xi) = 0$ for all $|\xi| > T$. Thus we have

$$(\mathfrak{M}_2 f)_{ij}(x) = \int_{\xi \in T_x M, \epsilon_2 < |\xi| < T} W_{ij}^{mn}(x, \xi) f_{mn}(\exp_x \xi) (1 - \chi_{\epsilon_2}(|\xi|)) d\mu_x(\xi).$$

It is clear that $\mathfrak{M}f = \mathfrak{M}_1 f + \mathfrak{M}_2 f$. This gives the promised decomposition of \mathfrak{M} . We next study \mathfrak{M}_1 and \mathfrak{M}_2 separately.

We first investigate \mathfrak{M}_1 . In the uniformly normal neighborhood of x_* , by a change of coordinate $y = \exp_x \xi$ and using the fact that $\frac{\xi}{|\xi|} = -\nabla_x \rho(x, y)$, $\exp_x \xi = \nabla_y \rho(x, y)$, we can deduce that (see [31]) \mathfrak{M}_1 has the following representation:

$$(\mathfrak{M}_1 f)_{ij}(x) = \frac{1}{\sqrt{\det g(x)}} \int_M K_{ij}^{mn}(x, y) f_{mn}(y) \chi_{\epsilon_1}(\rho(x, y)) dy, \quad (4.4)$$

where

$$K_{ij}^{mn}(x, y) = \frac{\partial \rho}{\partial x^i} \frac{\partial \rho}{\partial x^j} g^{mm'}(y) g^{nn'}(y) \frac{\partial \rho}{\partial x^{m'}} \frac{\partial \rho}{\partial x^{n'}} \cdot \left(\Phi^\dagger(x, -\nabla_x \rho) \Phi(y, \nabla_y \rho) + \Phi^\dagger(x, \nabla_x \rho) \Phi(x, -\nabla_y \rho) \right).$$

Following the same argument as in the proof of Lemma 3 in [33], we conclude that the following result holds.

Lemma 4.4. *\mathfrak{M}_1 is a classic ΨDO of order -1 in a neighborhood of x_* with principal symbol*

$$\sigma_p(\mathfrak{M}_1)_{ij}^{mn}(x, \omega) = \pi \int_{S_x M} (\Phi^\dagger \Phi)(x, \xi) \xi_i \xi_j \xi^m \xi^n \delta(\omega \cdot \xi) d\mu_x(\xi). \quad (4.5)$$

Here for each fixed i, j, m, n , $\sigma_p(\mathfrak{M}_1)_{ij}^{mn}(x, \omega)$ is a matrix from \mathbb{R}^{2d} to \mathbb{R}^{2d} .

We now proceed to study the operator \mathfrak{M}_2 whose property is determined by the exponential map $\exp(x_*, \cdot)$. We shall study the operator germ \mathfrak{M}_{2, ξ_*} for each $\xi \in T_{x_*} M$. We first consider the case when ξ_* is not a conjugate vector, i.e. ξ_* is a regular vector.

Lemma 4.5. *Let $\xi_* \in S_{x_*}^* \mathbb{R}^d$ be a regular vector, then there exists a neighborhood $U(x_*)$ of x_* and a neighborhood $B(x_*, \xi_*)$ of (x_*, ξ_*) such that for any $\chi \in C_0^\infty(B(x_*, \xi_*))$ the following operator*

$$(\mathfrak{M}_{2, \xi_*} f)_{ij}(x) = \int_{T_x M} W_{ij}^{mn}(x, \xi) f_{mn}(x, \xi) (1 - \chi_{\epsilon_2}(|\xi|)) \cdot \chi(x, \xi) d\mu_x(\xi)$$

is a smoothing operator from $\mathcal{E}'(S(\tau_2 M), \mathbb{R}^{2d})$ into $C^\infty(S(\tau_2 \overline{U(x_*)}), \mathbb{R}^{2d})$.

We next consider the case when ξ_* is a fold vector. We have the following result.

Lemma 4.6. *Let ξ_* be a fold vector of the map $\exp_{x_*}(\cdot)$. Then there exists a small neighborhood $U(x_*)$ of x_* and a small neighborhood $B(x_*, \xi_*)$ of (x_*, ξ_*) in \mathbb{R}^{2d} such that for any $\chi \in C_0^\infty(B(x_*, \xi_*))$, the operator $\mathfrak{M}_{2, \xi_*} : \mathcal{E}'(S(\tau_2 M), \mathbb{R}^{2d}) \rightarrow \mathcal{D}'(S(\tau_2 U(x_*)), \mathbb{R}^{2d})$ defined by*

$$(\mathfrak{M}_{2, \xi_*} f)_{ij}(x) = \int_{T_x M} W_{ij}^{mn}(x, \xi) f_{mn}(x, \xi) (1 - \chi_{\epsilon_2}(|\xi|)) \cdot \chi(x, \xi) d\mu_x(\xi), \quad f \in \mathcal{E}'(S(\tau_2 M), \mathbb{R}^{2d}) \quad (4.6)$$

is an FIO of order $-\frac{d}{2}$ whose associated canonical relation is compactly supported in the following set

$$\left\{ (x, \xi, y, \eta); \quad x \in U(x_*), y = \exp(x, \omega), (x, \omega) \in B(x_*, \xi_*), \det d_\omega \exp(x, \omega) = 0, \right. \\ \left. \xi = -\eta_i \frac{\partial \exp^i(x, \omega)}{\partial x}, \eta \in \text{Coker}(d_\omega \exp(x, \omega)). \right\} \quad (4.7)$$

Moreover, the canonical relation is the graph of a homogeneous canonical transformation from a neighborhood $(\exp_{x_} \xi_*, \exp_{x_*} \xi_*) \in T\mathbb{R}^d$ to $(x_*, \xi_*) \in T\mathbb{R}^d$, and hence \mathfrak{M}_{2, ξ_*} is bounded from H_{comp}^s to $H_{\text{local}}^{s+d/2}$.*

Proof. We note that for each fixed index quadruple (i, j, m, n) , the Schwartz kernel of the integral operator defined on the right hand side of (4.6) has the same type of singularities as those discussed in [35, 2]. Thus a similar argument as therein yields the result.

5 The X-ray transform operator \mathfrak{X} and its normal \mathfrak{N}

We study the X-ray transform \mathfrak{X} and its normal \mathfrak{N} in this section. Our goal is to prove Theorem 3.1 and Corollary 3.1 and 3.2.

We first show that the X-ray transform \mathfrak{X} vanishes on the potential fields.

Lemma 5.1. $\mathfrak{X}\mathcal{P} = 0$.

Proof. It suffices to show that $\mathfrak{X}d^s v = 0$ for all $v \in C_0^\infty(M)$. Indeed, let $v \in C_0^\infty(M)$. Consider the one-parameter family of diffeomorphisms $\phi^\tau : M \rightarrow M$ defined by

$$\frac{d\phi^\tau}{d\tau} = v, \quad \phi^0 = Id.$$

It is clear that $\phi^\tau|_{\partial M} = Id$. Denote $g^\tau = (\phi^\tau)^* g$ and $\mathcal{H}(\tau) = \mathcal{H}(g^\tau)$. We have $\mathcal{H}(\tau) \equiv \mathcal{H}(0)$. Thus $\mathcal{H}'(0) = 0$. On the other hand, a direct calculation shows that $\frac{dg^\tau}{d\tau}|_{\tau=0} = d^s v$, which further implies that $\mathcal{H}'(0) = \mathfrak{X} \frac{dg^\tau}{d\tau}|_{\tau=0} = \mathfrak{X}d^s v$. Therefore, we can conclude that $\mathfrak{X}d^s v = 0$, which completes the proof of the lemma.

Note that $\mathfrak{X} = \mathfrak{I} \circ \iota \circ \mathcal{L}$. We next present a decomposition of its normal \mathfrak{N} which follows from that of the operator \mathfrak{M} :

$$\mathfrak{N} = \mathfrak{X}^\dagger \mathfrak{X} = (\iota \circ \mathcal{L})^\dagger \circ \mathfrak{M} \circ (\iota \circ \mathcal{L}) = (\iota \circ \mathcal{L})^\dagger \circ \mathfrak{M}_1 \circ (\iota \circ \mathcal{L}) + (\iota \circ \mathcal{L})^\dagger \circ \mathfrak{M}_2 \circ (\iota \circ \mathcal{L}) = \mathfrak{N}_1 + \mathfrak{N}_2.$$

Similar to the truncation of the operator \mathfrak{X} , we introduce a truncation for the operator \mathfrak{J} . For any $\alpha \in C_0^\infty(S_- \partial M)$, we define

$$\mathfrak{J}_\alpha f(x_0, \xi_0) = \alpha(x_0, \xi_0) \mathfrak{J}_\alpha f(x_0, \xi_0). \quad (5.1)$$

It is clear that

$$\mathfrak{X}_\alpha = \mathfrak{J}_\alpha \circ \iota \circ \mathcal{L}.$$

By replacing the weight W_{ij}^{mn} with the new one $\alpha^\# \cdot W_{ij}^{mn}$, where $\alpha^\#$ is the unique lift of α to $S\mathbb{R}^d$ which is constant along each orbit of the geodesic flow $\mathcal{H}^t(g)$, we can define \mathfrak{M}_α , $\mathfrak{M}_{1,\alpha}$, $\mathfrak{M}_{2,\alpha}$ and consequently \mathfrak{N}_α , $\mathfrak{N}_{\alpha,1}$ and $\mathfrak{N}_{\alpha,2}$. It is clear that

$$\mathfrak{N}_\alpha = (\iota \circ \mathcal{L})^\dagger \circ \mathfrak{M}_\alpha \circ (\iota \circ \mathcal{L}) = (\iota \circ \mathcal{L})^\dagger \circ \mathfrak{M}_{\alpha,1} \circ (\iota \circ \mathcal{L}) + (\iota \circ \mathcal{L})^\dagger \circ \mathfrak{M}_{\alpha,2} \circ (\iota \circ \mathcal{L}) = \mathfrak{N}_{\alpha,1} + \mathfrak{N}_{\alpha,2}.$$

5.1 Local properties of the normal operator \mathfrak{N}

Let $x \in \tilde{M} \Subset M$ be a fold-regular point with the compact subset $\mathcal{Z}(x_*) \subset S_{x_*}M$ in Definition 3.2. We now construct a cut-off function $\alpha \in C_0^\infty(S_- \partial M)$ such that in a neighborhood of x_* , $\mathfrak{N}_{\alpha,1}$ is elliptic and $\mathfrak{N}_{\alpha,2}$ is smoother than $\mathfrak{N}_{\alpha,1}$ (see Lemma 5.5 and 5.6). The idea is to select a complete set of geodesics with no conjugate points except strong fold-regular ones. We follow closely the argument in [2].

Denote $\mathcal{C}_{\epsilon_2, T}\mathcal{Z} = \{r\xi; \xi \in \mathcal{Z}(x_*), r \in \mathbb{R} \text{ and } \epsilon_2 \leq |r| \leq T\}$. With the help of Lemma 4.5 and Lemma 4.6, there exist a finite number of vectors $\xi_j \in \mathcal{C}_{\epsilon_2, T}\mathcal{Z}$, $j = 1, 2, \dots, N$ such that for each ξ_j , there exist two neighborhoods $B_0(x_*, \xi_*) \Subset B(x_*, \xi_*)$ of $(x_*, \xi_j) \in \mathbb{R}^{2d}$ and a function $\chi_j \in C_0^\infty(B(x_*, \xi_j))$ with the following conditions:

(1). The operator \mathfrak{M}_{2, ξ_j} defined by

$$(\mathfrak{M}_{2, \xi_j} f)_{ij}(x) = \int_{T_x M} W_{ij}^{mn}(x, \xi) f_{mn}(\exp_x \xi) (1 - \chi_{\epsilon_2}(|\xi|)) \cdot \chi_j(x, \xi) d\mu_x(\xi)$$

is bounded from $H^k(S(\tau_2(M, M), \mathbb{R}^{2d}))$ to $H^{k+\frac{d}{2}}(S(\tau_2 U(x_*, \xi_*), \mathbb{R}^{2d}))$.

(2). $\mathcal{C}_{\epsilon_2, T}\mathcal{Z} \subset \bigcup_{j=1}^N B_0(x_*, \xi_j)$.

Denote by \mathcal{A}_0 be the greatest connected open symmetric subset in $\bigcup_{j=1}^N B_0(x_*, \xi_j)$ which contains $\mathcal{C}_{\epsilon_2, T}\mathcal{Z}$. Here and after, we say that a set \mathcal{B} in \mathbb{R}^{2d} is symmetric if $(x, \xi) \in \mathcal{B}$ implies that $(x, -\xi) \in \mathcal{B}$. Define

$$\mathcal{A}_\epsilon = \{(x, \xi) \in \mathbb{R}^{2d} : |x - x_*| \leq \epsilon, \epsilon_2 \leq |\xi| \leq T\}$$

for each $\epsilon > 0$. It is clear that \mathcal{A}_ϵ is compact in \mathbb{R}^{2d} , so is the set $\mathcal{A}_\epsilon \setminus \mathcal{A}_0$.

Lemma 5.2. *There exist $\epsilon_3 > 0$ and $\alpha \in C_0^\infty(S_- \partial M)$ such that the following two conditions are satisfied:*

$$\alpha(x_0, \xi_0) = 1 \quad \text{for all } (x_0, \xi_0) \in \mathcal{H}^t(g)(\mathcal{Z}(x_*)), \quad (5.2)$$

$$\alpha(x_0, \xi_0) = 0 \quad \text{for all } (x_0, \xi_0) \in \mathcal{H}^t(g)(\mathcal{A}_{\epsilon_3} \setminus \mathcal{A}_0). \quad (5.3)$$

Proof. See [2].

Lemma 5.3. *There exists a neighborhood $U(x_*)$ of x_* such that*

$$\ker \sigma_p(\mathfrak{M}_{\alpha,1})(x, \omega) = \{v_i^k \omega_j + v_j^k \omega_i : v_i^k \in \mathbb{R}, 1 \leq i, j \leq d, 1 \leq k \leq 2d\}, \quad (5.4)$$

for all $x \in U(x_*)$ and $\omega \in T_{x_*}^* M \setminus 0$.

Proof. Let $f \in \ker \sigma_p(\mathfrak{M}_{\alpha,1})(x, \omega)$, then

$$\langle f, \sigma_p(\mathfrak{M}_{\alpha,1})(x, \omega) f \rangle = 0. \quad (5.5)$$

Similar to Lemma 4.4, the symbol of $\mathfrak{M}_{\alpha,1}$ has the following form:

$$\sigma_p(\mathfrak{M}_{\alpha,1})^{mn}(x, \omega) = \pi \int_{\xi \in S_x M} |\alpha^\sharp(x, \xi)|^2 (\Phi^\dagger \circ \Phi)(x, \xi) \xi_i \xi_j \xi^m \xi^n \delta(\omega \cdot \xi) d\mu_x(\xi). \quad (5.6)$$

Substituting (5.6) into (5.5), we can deduce that

$$\int_{\xi \in S_x M, \xi \perp \omega} |\alpha^\sharp(x, \xi)|^2 \langle \Phi(x, \xi) f_{ij} \xi^i \xi^j, \Phi(x, \xi) f_{ij} \xi^i \xi^j \rangle d\mu_x(\xi) = 0.$$

We now show that (5.4) holds for the point x_* . Indeed, by Condition (5.2) and the “completeness” property of the set $\mathcal{Z}(x_*)$, there exists $\xi_0 \in S_{x_*} M$ with $\xi_0 \perp \omega$ such that $\alpha^\sharp(x, \xi_0) = 1$. It follows that

$$\langle \Phi(x, \xi) f_{ij} \xi^i \xi^j, \Phi(x, \xi) f_{ij} \xi^i \xi^j \rangle = 0$$

for ξ in a small neighborhood of ξ_0 such that $\xi \in S_x M$ and $\xi \perp \omega$. Since $\Phi(x, \xi)$ is invertible, this further implies that

$$f_{ij} \xi^i \xi^j = 0, \quad \text{for } \xi \text{ in a small neighborhood of } \xi_0 \text{ such that } \xi \in S_x M \text{ and } \xi \perp \omega.$$

Note that for each $k = 1, 2, \dots, 2d$, the following decomposition holds for $f^k = \{f_{ij}^k; 1 \leq i, j \leq d\}$:

$$f^k = h^k + i_\omega v^k, \quad \text{with } j_\omega h^k = 0.$$

Thus, $(h_{ij}^k + (i_\omega v^k)_{ij}) \xi^i \xi^j = 0$ for ξ in a small neighborhood of ξ_0 such that $\xi \in S_x M$ and $\xi \perp \omega$. This together with the fact that $(i_\omega v^k)_{ij} \xi^i \xi^j = 0$ yield $h_{ij}^k \xi^i \xi^j = 0$ for $\xi \perp \omega$. On the other hand, the equality $j_\omega h^k = 0$ implies that $h_{ij}^k \omega^i = 0$. Therefore, we can conclude that $h_{ij}^k \xi^i \xi^j = 0$ for all ξ in a small neighborhood of ξ_0 . It follows that $h^k = 0$. As a result, we get $f^k = i_\omega v^k = \frac{1}{2}(v_i^k \omega_j + v_j^k \omega_i)$. This proves the lemma for the point x_* .

To prove (5.4) for other points, we exploit the continuity of the function α^\sharp . In fact, we can find a neighborhood $U(x_*)$ of x_* such that for each $x \in U(x_*)$, there exists a “complete” set $\mathcal{Z}(x) \in S_x M$ such that $\alpha^\sharp(x, \xi) > \frac{1}{2}$ for all $\xi \in \mathcal{Z}(x)$. Then a similar argument as for the point x_* proves that (5.4) holds for these points in $U(x_*)$. This completes the proof of the lemma.

Lemma 5.4. $\mathfrak{N}_{\alpha,1}$ is a ΨDO of order one and there exists a neighborhood $U(x_*)$ of x_* such that for $x \in U(x_*)$,

$$\ker \sigma_p(\mathfrak{N}_{\alpha,1})(x, \omega) = \{f \in S_2 : f_{ij} = v_i \omega_j + v_j \omega_i \text{ for some } v \in S_1\}.$$

Moreover, $\sigma_p(\mathfrak{N}_{\alpha,1})(x, \omega)$ is positive definite on the set $\{f \in S_2 : f_{ij} \omega^i = 0\}$.

Proof. We first recall that $\mathfrak{N}_{\alpha,1} = (\iota \circ \mathcal{L})^\dagger \circ \mathfrak{M}_{\alpha,1} \circ (\iota \circ \mathcal{L})$. Since $\mathfrak{M}_{\alpha,1}$ is a Ψ DO of order -1 and both $\iota \circ \mathcal{L}$ and $(\iota \circ \mathcal{L})^\dagger$ are differential operators of order one, $\mathfrak{N}_{\alpha,1}$ is a Ψ DO of order one.

We next determine the kernel of $\sigma_p(\mathfrak{N}_{\alpha,1})$. Let $f = \{f_{ij}\} \in \ker \sigma_p(\mathfrak{N}_{\alpha,1})(x, \omega)$. Note that

$$\sigma_p(\mathfrak{N}_{\alpha,1}) = \sigma_p((\iota \circ \mathcal{L})^\dagger) \sigma_p(\mathfrak{M}_{\alpha,1}) \sigma_p(\iota \circ \mathcal{L}) = (\sigma_p(\iota \circ \mathcal{L}))^\dagger \sigma_p(\mathfrak{M}_{\alpha,1}) \sigma_p(\iota \circ \mathcal{L}).$$

We have

$$\sigma_p(\iota \circ \mathcal{L})(x, \omega) f \in \ker \sigma_p(\mathfrak{M}_{\alpha,1})(x, \omega).$$

By Lemma 5.3, there exists $v = \{v_i^k; 1 \leq k \leq 2d, 1 \leq i \leq d\}$ such that

$$(\sigma_p(\iota \circ \mathcal{L})(x, \omega) f)_{mn}^k = \frac{1}{2}(v_m^k \omega_n + v_n^k \omega_m), \quad k = 1, 2, \dots, 2d.$$

On the other hand, a direct calculation shows that

$$(\sigma_p(\iota \circ \mathcal{L})(x, \omega) f)_{mn}^{k+d} = \frac{1}{2} g^{kl} (f_{ml} \xi_n + f_{nl} \xi_m - f_{mn} \xi_l), \quad k = 1, 2, \dots, d.$$

Therefore,

$$f_{ml} \xi_n + f_{nl} \xi_m - f_{mn} \xi_l = v_{ml} \omega_n + v_{nl} \omega_m, \quad l = 1, 2, \dots, d. \quad (5.7)$$

where $v_{mn} := g_{ln}(x) v_m^{l+d}$, $1 \leq m, n, l \leq d$.

We now solve (5.7) for f . Let ω_a be the component of ω such that $\omega_a \neq 0$. First, by setting $m = n = a$ in (5.7), we obtain

$$2f_{al} \xi_a - f_{aa} \xi_l = 2v_{al} \omega_a. \quad (5.8)$$

Next, by setting $l = a$ in (5.8), we further get

$$f_{aa} = 2v_{aa}.$$

Substituting this back into (5.8), it follows that

$$f_{al} = \frac{\omega_l}{\omega_a} v_{aa} + v_{al}. \quad (5.9)$$

Finally, substituting (5.9) into (5.7) and letting $l = a$, we deduce that

$$\begin{aligned} f_{mn} &= \frac{2v_{aa}}{\omega_a^2} \omega_m \omega_n + \frac{\omega_m}{\omega_a} (v_{an} - v_{na}) + \frac{\omega_n}{\omega_a} (v_{am} - v_{ma}) \\ &= \frac{\omega_m}{\omega_a} (v_{an} - v_{na} + \frac{v_{aa}}{\omega_a} \omega_n) + \frac{\omega_n}{\omega_a} (v_{am} - v_{ma} + \frac{v_{aa}}{\omega_a} \omega_m) \\ &= \omega_m u_n + \omega_n u_m, \end{aligned}$$

where $u_m = \frac{1}{\omega_a} (v_{am} - v_{ma}) + \frac{v_{aa}}{\omega_a^2} \omega_m$. This gives the desired solution f . From this, we conclude that

$$\ker \sigma_p(\mathfrak{N}_{\alpha,1})(x, \omega) = \{f \in S_2 : f_{ij} = v_i \omega_j + v_j \omega_i\},$$

for $x \in U(x_*)$ where $U(x_*)$ is chosen as in Lemma 5.3. Finally, the remaining part of the lemma follows from the fact that $\sigma_p(\mathfrak{N}_{\alpha,1})(x, \omega)$ is symmetric and is non-negative. This completes the proof of the Lemma.

Lemma 5.5. *Let $x_* \in \tilde{M} \Subset M$ be a strong fold-regular point, $\alpha \in C_0^\infty(S_- \partial M)$ be chosen as in Lemma 5.2 and let k is a nonnegative integer. Then there exist two neighborhoods $U(x_*) \Subset \tilde{U}(x_*)$ of x_* such that the following estimate holds for all $f \in H^k(S(\tau_2(\tilde{M}, M)))$*

$$\|f^s\|_{H^k(U(x_*))} \lesssim \|\mathfrak{N}_{1,\alpha} f\|_{H^{k-1}(\tilde{U}(x_*))} + \|f\|_{H^{k-1}(M)}. \quad (5.10)$$

Proof. We divide the proof into the following three steps.

Step 1. We construct pseudo-inverse of the operator $\mathfrak{N}_{\alpha,1}$ for solenoidal 2-tensor fields. Consider the operator $\mathcal{A} = |D|^{-1} \mathfrak{N}_{\alpha,1} + \mathcal{P}$, where $|D|^{-1}$ is a properly supported parametrix for $(-\Delta_g)^{-\frac{1}{2}}$ in M . By Lemma 5.4, \mathcal{A} is an elliptic Ψ DO of order zero in a neighborhood $\tilde{U}(x_*)$ of x_* . We thus can find a Ψ DO of order zero, denoted by \mathcal{B} , such that

$$\mathcal{B}\mathcal{A} = \mathcal{Q} + \mathcal{K},$$

where \mathcal{Q} is a Ψ DO with symbol $\sigma(\mathcal{Q})(x, \omega) = Id$ for $x \in \tilde{U}(x_*)$ and \mathcal{K} is a smoothing operator. Let $\mathcal{K}_1 = \mathcal{K} + \mathcal{Q} - Id$. Then $\mathcal{B}\mathcal{A} = Id + \mathcal{K}_1$, i.e.,

$$\mathcal{B}|D|^{-1} \mathfrak{N}_{\alpha,1} + \mathcal{B}\mathcal{P} = Id + \mathcal{K}_1.$$

Applying \mathcal{S} from the right on both sides and using the fact that $\mathcal{P}\mathcal{S} = 0$, we obtain

$$\mathcal{B}|D|^{-1} \mathfrak{N}_{\alpha,1} \mathcal{S} = \mathcal{S} + \mathcal{K}_1 \mathcal{S}.$$

This completes the construction of the pseudo-inverse of $\mathfrak{N}_{\alpha,1}$.

Step 2. Let $f \in H^k(S(\tau_2(\tilde{M}, M)))$, we show that

$$\|f^s\|_{H^k(U(x_*))} \lesssim \|\mathfrak{N}_{1,\alpha} f^s\|_{H^{k-1}(\tilde{U}(x_*))} + \|f\|_{H^{-t}(M)} \quad (5.11)$$

for any $t > 0$. We argue as follows. By the result in Step 1, we have

$$f^s = \mathcal{S}f^s = \mathcal{B}|D|^{-1} \mathfrak{N}_{\alpha,1} f^s - \mathcal{K}_1 f^s \quad \text{in } \tilde{U}(x_*).$$

Let $U(x_*)$ be another neighborhood of x_* such that $U(x_*) \Subset \tilde{U}(x_*)$. Then

$$\|f^s\|_{H^k(U(x_*))} \leq \|\mathcal{B}|D|^{-1} \mathfrak{N}_{\alpha,1} f^s\|_{H^k(U(x_*))} + \|\mathcal{K}_1 f^s\|_{H^k(U(x_*))}.$$

We now estimate the term $\|\mathcal{B}|D|^{-1} \mathfrak{N}_{\alpha,1} f^s\|_{H^k(U(x_*))}$. Let $\chi_{U(x_*)}$ be a smooth function such that $\chi_{U(x_*)}|_{U_1(x_*)} = 1$ for some $U_1(x_*)$ satisfying $U(x_*) \Subset U_1(x_*) \Subset \tilde{U}(x_*)$, and $\chi_{U(x_*)} = 0$ outside $\tilde{U}(x_*)$. Note that

$$\mathcal{B}|D|^{-1} \mathfrak{N}_{\alpha,1} f^s = \mathcal{B}|D|^{-1} (\chi_{\tilde{U}(x_*)} \mathfrak{N}_{1,\alpha} f^s) + \mathcal{B}|D|^{-1} (1 - \chi_{\tilde{U}(x_*)}) \mathfrak{N}_{1,\alpha} f^s.$$

We can deduce that

$$\begin{aligned} \|\mathcal{B}|D|^{-1} (\chi_{U(x_*)} \mathfrak{N}_{1,\alpha} f^s)\|_{H^k(U(x_*))} &\lesssim \|\chi_{U(x_*)} \mathfrak{N}_{1,\alpha} f^s\|_{H^{k-1}(\tilde{U}(x_*))} \lesssim \|\mathfrak{N}_{1,\alpha} f^s\|_{H^{k-1}(\tilde{U}(x_*))}, \\ \|\mathcal{B}|D|^{-1} (1 - \chi_{U(x_*)}) \mathfrak{N}_{1,\alpha} f^s\|_{H^k(U(x_*))} &\lesssim \|f\|_{H^{-t}(M)}, \end{aligned}$$

where for the second inequality above we used the fact that the operator $\mathcal{B}|D|^{-1} (1 - \chi_{U(x_*)}) \mathfrak{N}_{1,\alpha} \mathcal{S}$

is smoothing from $H^k(S(\tau_2(\tilde{M}, M)))$ to $C^\infty(\overline{S(\tau_2\tilde{U}(x_*))})$. Thus, we conclude that

$$\|\mathcal{B}|D|^{-1}\mathfrak{N}_{\alpha,1}f^s\|_{H^k(U(x_*))} \lesssim \|\mathfrak{N}_{\alpha,1}f^s\|_{H^{k-1}(U(x_*))} + \|f\|_{H^{-t}(M)}.$$

for any $t > 0$.

To finish the proof of (5.11), it remains to show that $\mathcal{K}_1\mathcal{S}$ is a smoothing operator from $H^k(S(\tau_2(\tilde{M}, M)))$ to $C^\infty(S(\tau_2\tilde{U}(x_*)))$. Indeed, note that $\mathcal{K}_1 = \mathcal{K} + \mathcal{Q} - Id$. We need only show that $(\mathcal{Q} - Id)\mathcal{S}$ is smoothing from $L^2(S(\tau_2(\tilde{M}, M)))$ to $C^\infty(S(\tau_2\tilde{U}(x_*)))$. But this is a consequence of the fact that the symbol of $\mathcal{Q} - Id$ vanishes for $x \in \tilde{U}(x_*)$.

Step 3. We finally prove (5.10). Observe that $\mathfrak{N}_{\alpha,1}f^s = \mathfrak{N}_{\alpha,1}f - \mathfrak{N}_{\alpha,1}\mathcal{P}f$. By Lemma 5.4, the principle symbol of the operator $\mathfrak{N}_{\alpha,1}\mathcal{P}$ is equal to zero for $x \in \tilde{U}(x_*)$. As a result, $\mathfrak{N}_{\alpha,1}\mathcal{P}$ is a Ψ DO of order zero and therefore the estimate below holds

$$\|\mathfrak{N}_{\alpha,1}\mathcal{P}f\|_{H^{k-1}(U(x_*))} \lesssim \|f\|_{H^{k-1}(M)}.$$

It follows that

$$\|\mathfrak{N}_{\alpha,1}f^s\|_{H^{k-1}(U(x_*))} \lesssim \|\mathfrak{N}_{\alpha,1}f\|_{H^{k-1}(U(x_*))} + \|f\|_{H^{k-1}(M)}.$$

By substituting the above inequality into (5.11), we get (5.10). This completes the proof of the Lemma.

Lemma 5.6. *Let $x_* \in \tilde{M} \Subset M$ be a strong fold-regular point and $\alpha \in C_0^\infty(S_- \partial M)$ be chosen as in Lemma 5.2. Then there exists a neighborhood $U(x_*)$ of x_* such that the following estimate holds for all $f \in H^k(S^2(\tau_2(\tilde{M}, M)))$*

$$\|\mathfrak{N}_{2,\alpha}f\|_{H^{k-1}(U(x_*))} \lesssim \|f\|_{H^{k-\frac{d}{2}+1}(M)}$$

where k is any nonnegative integer.

Proof: With the help of Lemma 4.6, the proof is similar to that of Lemma 6.7 in [2].

5.2 Proof of Theorem 3.1 and Corollary 3.1 and 3.2

Proof of Theorem 3.1: It is a direct consequence of Lemma 5.5 and 5.6.

Proof of Corollary 3.1: For each $x \in \tilde{M}$, by Theorem 3.1, there exist neighborhoods $U(x) \Subset \tilde{U}(x_*)$ of x and a smooth function $\alpha \in C_0^\infty(S_- \partial M)$ such that the estimate (3.3) holds. Since \tilde{M} is compact, we can find finite number of points, say x_1, x_2, \dots, x_N , such that $\tilde{M} \subset \bigcup_{j=1}^N U(x_j)$ and the following estimate holds for each j :

$$\|f^s\|_{H^k(U(x_j))} \lesssim \|\mathfrak{N}_{\alpha_j}f\|_{H^{k-1}(\tilde{U}(x_j))} + \|f\|_{H^{k+1-\frac{d}{2}}(M)}.$$

Let $\tilde{M}_1 = \bigcup_{j=1}^N U(x_j)$. Then $\tilde{M} \Subset \tilde{M}_1$. We claim that

$$\|f^s\|_{H^k(M \setminus \tilde{M}_1)} \lesssim \|f\|_{H^{k+1-\frac{d}{2}}(M)}.$$

Indeed, note that $f^s = \mathcal{S}f = (Id - d^s(\Delta_M^s)^{-1}\delta^s)f$, where Δ_M^s is the Dirichlet realization of the operator $\Delta^s := \delta^s d^s$ in M and $(\Delta_M^s)^{-1}$ is its inverse. Using the pseudo-local property of the operator \mathcal{S} , we see that \mathcal{S} is smoothing from $L^2(\tilde{M})$ to $C^\infty(M \setminus \tilde{M}_1)$. Hence the claim follows.

Therefore, we deduce that

$$\begin{aligned}\|f^s\|_{H^k(M)} &\leq \|f^s\|_{H^k(M \setminus \tilde{M}_1)} + \sum_{j=1}^N \|f\|_{H^k(U(x_j))} \\ &\lesssim \sum_{j=1}^N \|\mathfrak{N}_{\alpha_j} f\|_{H^k(\tilde{U}(x_j))} + \|f\|_{H^{k+1-\frac{d}{2}}(M)}.\end{aligned}$$

The corollary is proved.

Proof of Corollary 3.2:

Step 1. Let x_j , $U(x_j)$ and α_j be chosen as in Theorem 3.1. Denote by H the Hilbert space $\prod_{j=1}^N H^{-1}(S(\tau_2 \tilde{U}(x_j)))$. We consider the operator

$$T : L^2(S(\tau_2(\tilde{M}, M))) \rightarrow H$$

defined by

$$Tf = (\mathfrak{N}_{\alpha_1} f, \mathfrak{N}_{\alpha_2} f, \dots, \mathfrak{N}_{\alpha_N} f).$$

Then Theorem 3.1 can be interpreted as

$$\|f^s\|_{L^2(M)} \lesssim \|Tf\|_H + \|f\|_{H^{1-\frac{d}{2}}(M)}. \quad (5.12)$$

Note that the operator T vanishes on the potential vector fields (see Lemma 5.1). We can define $T_0 : \mathcal{SL}^2(S(\tau_2(\tilde{M}, M))) \rightarrow H$ by

$$T_0 f^s = Tf.$$

Step 2. We show that the space $\mathfrak{L} := \ker T_0$ is finite dimensional in $\mathcal{SL}^2(S(\tau_2(\tilde{M}, M)))$. Indeed, by contradiction, assume that there exist a pairwise orthogonal sequence $f_n^s \in \mathfrak{L}$ such that $\|f_n^s\|_{L^2(M)} = 1$ and $T_0 f_n^s = 0$. By Lemma 5.7, there exist $C > 0$, $g_n \in L^2(S(\tau_2(\tilde{M}, M)))$ such that

$$\|g_n\|_{L^2(M)} \leq C, \quad \text{and} \quad g_n^s = f_n^s.$$

Then the inequality (5.12) yields that

$$\|f_n^s - f_m^s\|_{L^2(M)} \lesssim \|g_n - g_m\|_{H^{1-\frac{d}{2}}(M)}.$$

Since for $d \geq 3$, $H^{1-\frac{d}{2}}(M)$ is compactly embedded into $L^2(M)$, we can find a subsequence of g_n , still denoted by g_n , such that g_n is Cauchy in $H^{1-\frac{d}{2}}(M)$. Then f_n^s is also Cauchy in $L^2(M)$. This contradicts to the assumption that f_n^s are pairwise orthogonal in $L^2(S(\tau_2(\tilde{M}, M)))$. This contradiction proves our assertion.

Step 3. We finally show (3.4). Assume the contrary, then there exists a sequence $f_n \in L^2(S(\tau_2(\tilde{M}, M)))$ such that $\{f_n^s\}_{n=1}^\infty \subset \mathcal{SL}^2(S(\tau_2(\tilde{M}, M))) \cap \mathfrak{L}^\perp$, $\|f_n^s\|_{L^2(M)} = 1$ and $\|Tf_n\|_H \leq \frac{1}{n}$ for all n . By the same argument as in Step 2, we may assume that the sequence $\{f_n\}_{n=1}^\infty$ is Cauchy in $L^2(S\tau_2 M)$. As a result, we can conclude that $\{f_n^s\}_{n=1}^\infty$ is Cauchy in $L^2(M)$ by using (5.12). Note that $\mathcal{SL}^2(S(\tau_2(\tilde{M}, M)))$ is closed in $L^2(S\tau_2 M)$ by Lemma 5.8. There exists

$f_0 \in L^2(S(\tau_2(\tilde{M}, M)))$ such that $f_0^s = \lim_{n \rightarrow \infty} f_n^s$. It is clear that $f_0^s \in \mathfrak{L}$. However, since \mathfrak{L}^\perp is closed, as the limit of a sequence of functions in \mathfrak{L}^\perp , f_0^s should also belong to \mathfrak{L}_0^\perp . Therefore, $f_0^s = 0$. But this contradicts to the fact that $\|f_0^s\|_{L^2(M)} = \lim_{n \rightarrow \infty} \|f_n^s\|_{L^2(M)} = 1$. This contradiction completes the proof of (3.4) and hence the corollary.

Proof of Corollary 3.3: The same argument as for Corollary 3.2.

We now present two axillary lemmas which are needed in the proofs above.

Lemma 5.7. *Let k be a nonnegative integer. Then there exists $C > 0$ such that for any $f^s \in \mathcal{SH}^k(S(\tau_2(\tilde{M}, M)))$, we can find $g \in H^k(S(\tau_2(\tilde{M}, M)))$ such that*

$$\|g\|_{H^k(M)} \leq C \|f^s\|_{H^k(M)}, \quad \text{and } g^s = f^s.$$

Proof. We first prove the lemma for the case $k = 0$. Let $f = f^s + d^s v$ where $f \in L^2(S\tau_2(\tilde{M}, M))$ and $v \in H_0^1(S(\tau_2(\tilde{M}, M)))$. We need find $g \in L^2(S(\tau_2(\tilde{M}, M)))$ such that $g^s = f^s$ and $\|g\|_{L^2(M)} \leq C \|f^s\|_{L^2(M)}$. We argue as follows.

We first claim that

$$\|v\|_{H^1(M \setminus \tilde{M})} \lesssim \|f^s\|_{L^2(M)}. \quad (5.13)$$

Indeed, noting that $f = 0$ in $M \setminus \tilde{M}$, we have

$$d^s v = -f^s \quad \text{in } M \setminus \tilde{M}.$$

Using $v|_{\partial M} = 0$ and the same argument as in [31], we can show that

$$\|v\|_{L^2(M \setminus \tilde{M})} \lesssim \|f^s\|_{L^2(M \setminus \tilde{M})}. \quad (5.14)$$

On the other hand, observe that

$$\begin{aligned} (d^s v)_{ij} &= \frac{1}{2} (\nabla_i v_j + \nabla_j v_i) \\ &= \frac{1}{2} \left(\frac{\partial v_j}{\partial x^i} + \frac{\partial v_i}{\partial x^j} - 2\Gamma_{ij}^k v_k \right). \end{aligned}$$

Thus,

$$\frac{\partial v_j}{\partial x^i} + \frac{\partial v_i}{\partial x^j} = 2(d^s v)_{ij} + 2\Gamma_{ij}^k v_k = -2f^s + 2\Gamma_{ij}^k v_k \quad \text{in } M \setminus \tilde{M}.$$

Consequently,

$$\left\| \frac{\partial v_j}{\partial x^i} + \frac{\partial v_i}{\partial x^j} \right\|_{L^2(M \setminus \tilde{M})} \lesssim \|f^s\|_{L^2(M \setminus \tilde{M})} + \|v\|_{L^2(M \setminus \tilde{M})} \lesssim \|f^s\|_{L^2(M \setminus \tilde{M})}. \quad (5.15)$$

By using the following Korn's inequality (see for instance [14])

$$\|v\|_{H^1(M \setminus \tilde{M})} \lesssim \|v\|_{L^2(M \setminus \tilde{M})} + \sum_{ij} \left\| \frac{\partial v_j}{\partial x^i} + \frac{\partial v_i}{\partial x^j} \right\|_{L^2(M \setminus \tilde{M})}, \quad (5.16)$$

the estimate (5.13) follows immediately from (5.14)-(5.16). This completes the proof of the claim.

Next, for each $v \in H^1(M \setminus \tilde{M})$, we can find an extension $\tilde{v} \in H^1(M)$ such that $\tilde{v} = v$ in

$M \setminus \tilde{M}$ and $\|\tilde{v}\|_{H^1(M)} \lesssim \|v\|_{H^1(M \setminus \tilde{M})}$. Let $g = f^s + \tilde{v}$. Then $g \in L^2(S(\tau_2(\tilde{M}, M)))$ and $g^s = f^s$. Moreover, $\|g\|_{L^2(M)} \lesssim \|f^s\|_{L^2(M)}$. The lemma is proved for the case $k = 0$.

Finally, for the case $k \geq 1$, the argument is almost the same. The key point is to derive the following estimate

$$\|v\|_{H^{k+1}(M \setminus \tilde{M})} \lesssim \|d^s v\|_{H^k(M \setminus \tilde{M})},$$

by Korn's inequality and the estimate (5.14).

As a direct consequence of the above lemma, we obtain the following result.

Lemma 5.8. *The linear space $\mathcal{S}H^k(S(\tau_2(\tilde{M}, M)))$ is closed in $H^k(S(\tau_2 M))$.*

6 Proof of Theorem 1.1

We write $\tilde{g} = g + f$. Then $\text{supp } f \subset \hat{M}$. We assume that f is sufficiently small in $C^{k,s}$ norm for some sufficiently large integer k and $0 < s < 1$ (we will choose $k = 15$ and $s = \frac{1}{2}$ at the end of the proof). We divide the argument into the following eight steps.

Step 1. We first construct a diffeomorphism ϕ of M such that $\phi^* \tilde{g}$ is solenoidal with respect to the metric g . In fact, by solving the equation $\delta^s(\phi^* \tilde{g} - g) = 0$, we obtain a unique ϕ which is close to the identity map and equal to it on the boundary ∂M and satisfies the following estimate

$$\|\phi - Id\|_{C^{k,s}} \lesssim \|f\|_{C^{k,s}}. \quad (6.1)$$

We denote $\tilde{g}_1 = \phi^* \tilde{g}$ and $f_1 = \tilde{g}_1 - g$. Then the following estimate holds

$$\|f_1\|_{C^{k-1,s}} \lesssim \|f\|_{C^{k,s}}. \quad (6.2)$$

Step 2. We construct another diffeomorphism ψ of M such that $\psi^* \tilde{g}_1 = e$ in a small neighborhood of ∂M . We follow the same approach as in Section 4.2 in [34]. Denote by $\exp_{\partial M, g}$ and $\exp_{\partial M, \tilde{g}_1}$ the boundary normal coordinate near ∂M with respect to the metric g and \tilde{g}_1 , respectively. Both maps are well-defined from $\partial M \times [0, \epsilon_0]$ for some ϵ_0 , which can be chosen to be independent of \tilde{g} , to some neighborhoods of ∂M . Let $\psi_1 = \exp_{\partial M, g}(\exp_{\partial M, \tilde{g}_1})^{-1}$. Then ψ_1 maps a small neighborhood of ∂M to another and satisfies the following estimate:

$$\|\psi_1 - Id\|_{C^{k-2}} \lesssim \|f_1\|_{C^{k-1}}. \quad (6.3)$$

If $\|f_1\|_{C^{k-1}}$ is sufficiently small, we can show that there exists a unique vector field $v \in C^{k-2}$ in a sufficiently small neighborhood of ∂M , say W_1 , with $v = 0$ on ∂M such that

$$\psi_1(x) = \exp_x v(x)$$

and

$$\|v\|_{C^{k-2}} \lesssim \|\psi_1 - Id\|_{C^{k-2}} \lesssim \|f_1\|_{C^{k-1}}.$$

We now extend ψ_1 to a diffeomorphism of the whole domain M . Let W be another neighborhood of ∂M such that $W \Subset W_1 \cap (M \setminus \tilde{M})$. We may choose W such that $M \setminus \overline{W}$ is a strictly convex and smooth sub-domain of M . Let χ be fixed a smooth cutoff function such that $\chi = 1$

on W and $\chi = 0$ outside W_1 . We define

$$\psi(x) = \exp_x(\chi(x) \cdot v(x)).$$

Then ψ is a diffeomorphism which equals to the identity map on the boundary and satisfies the following estimate

$$\|\psi - Id\|_{C^{k-2}} \lesssim \|v\|_{C^{k-2}} \lesssim \|f_1\|_{C^{k-1}}. \quad (6.4)$$

This finishes the construction of ψ . We note that the diffeomorphism ψ constructed above depends only on \tilde{g} and the choice of the cutoff function χ .

Step 3. Denote $\tilde{g}_2 = \psi^* \tilde{g}_1$. It is clear that \tilde{g}_2 is isometric to \tilde{g}_1 and hence \tilde{g} . Moreover, \tilde{g}_2 has the same boundary normal coordinate as g and hence \tilde{g} also in W (both g and \tilde{g} equal to e there). Therefore, we can conclude that $\tilde{g}_2 = \tilde{g} = e$ in W .

Let $\tilde{M} = M \setminus W$ and $f_2 = \tilde{g}_2 - g$. Then the support of f_2 is contained in \tilde{M} .

We define $\varphi = \phi \circ \psi$. It is clear that $\tilde{g}_2 = \psi^* \tilde{g}_1 = \psi^* \phi^* \tilde{g} = \varphi^* \tilde{g}$.

Step 4. We now have constructed two isometric copies of \tilde{g} such that one is solenoidal with respect to the metric g and the other is equal to e in the neighborhood W of ∂M . Moreover, the following estimate holds for $f_1 = \tilde{g}_1 - g$ and $f_2 = \tilde{g}_2 - g$:

$$\|f_2 - f_1\|_{C^{l-3}} = \|\tilde{g}_2 - \tilde{g}_1\|_{C^{l-3}} \lesssim \|\psi - Id\|_{C^{l-2}} \lesssim \|f_1\|_{C^{l-1}}, \quad \text{for all } 3 \leq l \leq k. \quad (6.5)$$

By Proposition 1 in [34], we also have

$$\|f_2^s - f_1\|_{C^{l-3}} \lesssim \|f_1\|_{C^{l-1}}^2, \quad \text{for all } 3 \leq l \leq k. \quad (6.6)$$

Step 5. Note that \tilde{g}_2 equals to e in the neighborhood W of ∂M . By Proposition 1 and the fact that $\mathcal{H}^T(\tilde{g}_2) = \mathcal{H}^T(\tilde{g})$, we obtain

$$\mathcal{H}^T(\tilde{g}_2) - \mathcal{H}^T(g) = \mathcal{H}^T(\tilde{g}) - \mathcal{H}^T(g) = \mathfrak{X}f_2 + O(\|f_2\|_{C^{5+[\frac{d}{2}]}(M)}^2) \quad \text{in } C^{3+[\frac{d}{2}]}(S_- \partial M). \quad (6.7)$$

Step 6. With the sub-manifold \tilde{M} defined in Step 3, we let $\tilde{U}(x_j) \subset M$, $\alpha_j \in C_0^\infty(S_- \partial M)$, $j = 1, 2, \dots, N$, and the finite dimensional space $\mathfrak{L} \in SL^2(S(\tau_2(M, \tilde{M})))$ be determined as in Corollary 3.2. By our assumption, $f_2 \perp \mathfrak{L}$, which further yields $f_2^s \perp \mathfrak{L}$ by using the orthogonal decomposition of f_2 . Therefore, the following estimate holds by Corollary 3.3:

$$\|f_2^s\|_{H^{3+[\frac{d}{2}]}(M)} \lesssim \sum_{j=1}^N \|\mathfrak{N}_{\alpha_j} f_2\|_{H^{2+[\frac{d}{2}]}(\tilde{U}(x_j))}. \quad (6.8)$$

Step 7. For each $j = 1, 2, \dots, N$, we have

$$\begin{aligned} \|\mathfrak{N}_{\alpha_j} f_2\|_{H^{2+[\frac{d}{2}]}(\tilde{U}(x_j))} &= \|\mathcal{L}^\dagger \mathfrak{I}_{\alpha_j}^\dagger \mathfrak{X}_{\alpha_j} f_2\|_{H^{2+[\frac{d}{2}]}(\tilde{U}(x_j))} \\ &\lesssim \|\mathfrak{I}_{\alpha_j}^\dagger \mathfrak{X}_{\alpha_j} f_2\|_{H^{3+[\frac{d}{2}]}(\tilde{U}(x_j))} \\ &\lesssim \|\mathfrak{I}_{\alpha_j}^\dagger \mathfrak{X}_{\alpha_j} f_2\|_{H^{3+[\frac{d}{2}]}(M)} \\ &\lesssim \|\mathfrak{X}_{\alpha_j} f_2\|_{H^{3+[\frac{d}{2}]}(S_- \partial M)}. \end{aligned} \quad (6.9)$$

Step 8. By (6.7), (6.8) and (6.9), we deduce that

$$\begin{aligned}\|f_2^s\|_{H^{3+[\frac{d}{2}]}(M)} &\lesssim \sum_{j=1}^N \|\mathfrak{X}_{\alpha_j} f_2\|_{H^{3+[\frac{d}{2}]}(S_{-\partial M})} \\ &\lesssim \sum_{j=1}^N \|\alpha_j(\mathcal{H}^T(\tilde{g}) - \mathcal{H}^T(g))\|_{H^{3+[\frac{d}{2}]}(S_{-\partial M})} + \|f_2\|_{C^{5+[\frac{d}{2}]}(M)}^2.\end{aligned}$$

Therefore,

$$\begin{aligned}\|f_1\|_{C^2(M)} &\lesssim \|f_2^s\|_{C^2(M)} + \|f_1\|_{C^4(M)}^2 \quad (\text{by taking } l = 5 \text{ in (6.6)}) \\ &\lesssim \|f_2^s\|_{H^{3+[\frac{d}{2}]}(M)} + \|f_1\|_{C^4(M)}^2 \quad (\text{by Sobolev embedding theorem}) \\ &\lesssim \sum_{j=1}^N \|\alpha_j(\mathcal{H}^T(\tilde{g}) - \mathcal{H}^T(g))\|_{H^{3+[\frac{d}{2}]}(S_{-\partial M})} + \|f_2\|_{C^{5+[\frac{d}{2}]}(M)}^2 + \|f_1\|_{C^4(M)}^2 \\ &\lesssim \sum_{j=1}^N \|\alpha_j(\mathcal{H}^T(\tilde{g}) - \mathcal{H}^T(g))\|_{H^{3+[\frac{d}{2}]}(S_{-\partial M})} + \|f_1\|_{C^{7+[\frac{d}{2}]}(M)}^2 \quad (\text{by taking } l = 8 + [\frac{d}{2}] \text{ in (6.5)}) \\ &\lesssim \sum_{j=1}^N \|\alpha_j(\mathcal{H}^T(\tilde{g}) - \mathcal{H}^T(g))\|_{H^{3+[\frac{d}{2}]}(S_{-\partial M})} + \|f_1\|_{C^2(M)} \cdot \|f_1\|_{C^{12+2[\frac{d}{2}]}(M)}.\end{aligned}$$

Note that $\|f_1\|_{C^{12+2[\frac{d}{2}]}(M)} \leq \|f_1\|_{C^{12+2[\frac{d}{2}], \frac{1}{2}}(M)} \lesssim \|f\|_{C^{13+2[\frac{d}{2}], \frac{1}{2}}(M)}$ (by taking $k = 13 + 2[\frac{d}{2}]$ in (6.2) for the last inequality). By letting $\|f\|_{C^{13+2[\frac{d}{2}], \frac{1}{2}}(M)}$ sufficiently small, we obtain

$$\|f_1\|_{C^2(M)} \lesssim \sum_{j=1}^N \|\alpha_j(\mathcal{H}^T(\tilde{g}) - \mathcal{H}^T(g))\|_{H^{3+[\frac{d}{2}]}(S_{-\partial M})}.$$

Finally, by taking $l = 3$ in (6.5), we get the desired estimate

$$\|f_2\|_{C^0(M)} \lesssim \|f_1\|_{C^2(M)} \lesssim \sum_{j=1}^N \|\alpha_j(\mathcal{H}^T(\tilde{g}) - \mathcal{H}^T(g))\|_{H^{3+[\frac{d}{2}]}(S_{-\partial M})},$$

which completes the proof of the theorem.

Appendix A: Relations between the geodesic flow and the scattering relation

We present the relations between the geodesic flow and the scattering relation in the form (1.1). Let \mathcal{H}_0^t denote the geodesic flow induced by the Euclidean metric in \mathbb{R}^d and let g be the background metric. Then the following identity holds:

$$\mathcal{H}^T(g)(x_0, \xi_0) = \mathcal{H}_0^{T-L(g)(x_0, \xi_0)}(\Sigma(g)(x_0, \xi_0)), \quad (x_0, \xi_0) \in S_{-\partial M}. \quad (6.10)$$

On the other hand, for each $(x, \xi) \in S(\mathbb{R}^d \setminus M)$, we define $\kappa(x, \xi)$ to be the first positive moment when the orbit of the geodesic flow $\mathcal{H}_0^t(x, -\xi)$ hits the boundary ∂M . By the assumptions on the background metric g , κ is well-defined and is smooth in a sufficiently small neighborhood of $\mathcal{H}^T(g)(x_0, \xi_0)$ for each $(x_0, \xi_0) \in S_- \partial M$. As a result, the following identities hold for all $(x_0, \xi_0) \in S_- \partial M$:

$$L(g)(x_0, \xi_0) = \kappa(\mathcal{H}^T(g)(x_0, \xi_0)); \quad (6.11)$$

$$\Sigma(g)(x_0, \xi_0) = \mathcal{H}_0^{\kappa(\mathcal{H}^T(g)(x_0, \xi_0))}(\mathcal{H}^T(g)(x_0, \xi_0)). \quad (6.12)$$

Now, by taking derivatives with respect to g in the above identities (6.10), (6.11) and (6.12), we obtain the following result.

Lemma 6.1. *There exist two smooth matrix-valued functions $B \in C^\infty(S_- \partial M)$ and $C \in C^\infty(S_- \partial M)$ such that for all $(x_0, \xi_0) \in S_- \partial M$,*

$$\frac{\delta \mathcal{H}^T(g)}{\delta g}(x_0, \xi_0) = B(x_0, \xi_0) \begin{pmatrix} \frac{\delta \Sigma(g)}{\delta g}(x_0, \xi_0) \\ \frac{\delta L(g)}{\delta g}(x_0, \xi_0) \end{pmatrix}; \quad \begin{pmatrix} \frac{\delta \Sigma(g)}{\delta g}(x_0, \xi_0) \\ \frac{\delta L(g)}{\delta g}(x_0, \xi_0) \end{pmatrix} = C(x_0, \xi_0) \frac{\delta \mathcal{H}^T(g)}{\delta g}(x_0, \xi_0).$$

Remark 6.1. *The matrix-valued function C is only smooth in $S_- \partial M$ and it may not be extended smoothly to $\overline{S_- \partial M}$. This is because the function κ may not have a smooth extension to $\mathcal{H}^T(g)(\overline{S_- \partial M})$.*

Appendix B: Proof of Proposition 1

We first linearize the nonlinear operator which maps a metric to its corresponding Christoffel symbol at g . We denote by $\Gamma(g + f)$ the Christoffel symbol corresponding to the metric $g + f$. Observe that

$$\Gamma_{ij}^k(g) = \frac{1}{2} \left(\frac{\partial g_{jp}}{\partial x^i} + \frac{\partial g_{ip}}{\partial x^j} - \frac{\partial g_{ij}}{\partial x^p} \right). \quad (6.13)$$

We can show that

$$\mathcal{L}f = \frac{\delta \Gamma(g)}{\delta g} f.$$

We now present an estimate for the quality of the linear approximation.

Lemma 6.2. *For $f \in S(\tau_2 M)$ with sufficiently small C^2 norm, the following estimate holds for any positive integer k .*

$$\|\Gamma(g + f) - \Gamma(g) - \mathcal{L}f\|_{C^k(M)} \lesssim \|f\|_{C^{k+1}(M)}^2.$$

Proof. For fixed g and f , define $\Gamma(\tau) = \Gamma(g + \tau f)$. It is clear that

$$\Gamma(g + f) - \Gamma(g) = \Gamma(1) - \Gamma(0) = \Gamma'(0) + \frac{1}{2} \Gamma''(\eta), \quad \text{for some } \eta \in [0, 1]$$

and $\Gamma'(0) = \mathcal{L}f$. To prove the lemma, it suffices to show that

$$\|\Gamma''(\tau)\|_{C^k} \lesssim \|f\|_{C^{k+1}}^2.$$

We argue as follows. First, from Formula (6.13), we can deduce that

$$2\Gamma_{ij}^k(\tau)g_{kl}(\tau) = \frac{\partial g_{jl}}{\partial x^i} + \frac{\partial g_{il}}{\partial x^j} - \frac{\partial g_{ij}}{\partial x^l} + \tau\left(\frac{\partial f_{jl}}{\partial x^i} + \frac{\partial f_{il}}{\partial x^j} - \frac{\partial f_{ij}}{\partial x^l}\right).$$

Next, we differentiate the above identity with respect to τ to obtain

$$2\dot{\Gamma}_{ij}^k(\tau)g_{kl}(\tau) + 2\Gamma_{ij}^k(\tau)f_{kl} = \frac{\partial f_{jl}}{\partial x^i} + \frac{\partial f_{il}}{\partial x^j} - \frac{\partial f_{ij}}{\partial x^l}. \quad (6.14)$$

We then multiply both sides of the above identity by g^{lp} and summate over l . This gives

$$2\dot{\Gamma}_{ij}^k(\tau) = -2\Gamma_{ij}^k(\tau)f_{kl}g^{lp} + g^{lp}\left(\frac{\partial f_{jl}}{\partial x^i} + \frac{\partial f_{il}}{\partial x^j} - \frac{\partial f_{ij}}{\partial x^l}\right).$$

It follows that

$$\|\dot{\Gamma}_{ij}^k(\tau)\|_{C^k} \lesssim \|f\|_{C^{k+1}}.$$

Now, differentiating the identity (6.14) with respect to τ again,

$$\ddot{\Gamma}_{ij}^k(\tau)g_{kl}(\tau) + 2\dot{\Gamma}_{ij}^k(\tau)f_{kl} = 0.$$

Using the same trick as we did for $\dot{\Gamma}_{ij}^k(\tau)$, we can conclude that

$$\|\ddot{\Gamma}_{ij}^k(\tau)\|_{C^k} \lesssim \|\dot{\Gamma}_{ij}^k(\tau)\|_{C^k}\|f\|_{C^k} \lesssim \|f\|_{C^{k+1}}^2,$$

which yields the desired estimate and completes the proof of the lemma.

We next linearize the nonlinear operator which maps a Christoffel symbol to its corresponding geodesic flow at Γ . By a similar argument as in [2], we can show that the following holds.

Lemma 6.3. *Let k be a nonnegative integer, the following estimate holds for the linearization of the nonlinear operator which maps Christoffel symbols to its corresponding geodesic flow at Γ :*

$$\|\mathcal{H}^T(\Gamma + \tilde{\Gamma}) - \mathcal{H}^T(\Gamma) - \mathfrak{J}(\tilde{\Gamma})\|_{C^k(S_{-\partial M})} \lesssim \|\tilde{\Gamma}\|_{C^{k+1}(M)}^2. \quad (6.15)$$

Proof of Proposition 1: It is a consequence of Lemma 6.2, 6.3 and the chain rule.

References

- [1] V. I. Arnold, S. M. Gusein-Zade and A. N. Varchenko. Singularities of Differentiable Maps. I, Volume 82 of Monographs in Mathematics. Birkhäuser, Boston, 1985.
- [2] G. Bao and H. Zhang. Sensitivity analysis of an inverse problem for the wave equation with caustics, Journal of AMS, accepted.
- [3] I. N. Bernstein and M. L. Gerver. Conditions of distinguishability of metrics by godographs, Methods and Algorithms of Interpretation of Seismological Information, Computerized Seismology, 13(1980), 50-73.
- [4] G. Besson, G. Courtois and S. Gallot. Entropies et rigidités des espaces localement symétriques de courbure strictement négative, Geom. Funct. Anal., 5(1995), 731-799.

- [5] G. Beylkin. Stability and uniqueness of the solution of the inverse kinematic problem in the multidimensional case, *J. Soviet Math.*, 21(1983), 251-254.
- [6] D. Burago and S. Ivanov. Boundary rigidity and filling volume minimality of metrics close to a flat one, *Ann. of Math.* 171(2)(2010), 1183-1211.
- [7] C. Croke. Rigidity for surfaces of non-positive curvature, *Comment. Math. Helv.*, 65(1990), 150-169.
- [8] C. Croke. Rigidity and the Distance Between Boundary Points, *J. Diff. Geom.*, 33(1991), 445-464.
- [9] C. Croke. Boundary and lens rigidity of finite quotients, *Proc. AMS*, 133(12)(2005), 3663-3668.
- [10] C. Croke. Scattering rigidity with trapped geodesics, *Ergod. Th. and Dynam. Sys*, 2013.
- [11] C. Croke, N. Dairbekov and V. Sharafutdinov. Local boundary rigidity of a compact Riemannian manifold with curvature bounded above, *Trans. AMS*, 352(9)(2000), 3937-3956.
- [12] C. Croke and P. Herreros. Lens rigidity with trapped geodesics in two dimensions, *arXiv: 1108.4938*, 2011.
- [13] C. Croke and B. Kleiner. Conjugacy and rigidity for manifolds with a parallel vector field, *J. Diff. Geom.*, 39(1994), 659-680.
- [14] G. Duvaut and G. L. Lions. *Inequalities in Mechanics and Physics*, Springer-Verlag, 1976.
- [15] G. Eskin. Inverse scattering problem in anisotropic media, *Comm. Math. Phys.*, 199(2)(1998), 471-491.
- [16] P. Gabriel, M. Salo and G. Uhlmann. Tensor tomography on simple surfaces, *Invent. Math.*, 193(2013), 229-247.
- [17] M. Gromov. Filling Riemannian manifolds, *J. Diff. Geom.*, 18(1)(1983), 1-148.
- [18] V. Guillemin. Sojourn times and asymptotic properties of the scattering matrix, *Proceedings of the Oji Seminar on Algebraic Analysis and the RIMS Symposium on Algebraic Analysis (Kyoto Univ., Kyoto, 1976)*. Publ. Res. Inst. Math. Sci. 12(1976/77), supplement, 69-88.
- [19] L. Hörmander. *The Analysis of Linear Partial Differential Operators*, III and IV, Springer-Verlag, Berlin, 1985.
- [20] M. Lassas, V. Sharafutdinov and G. Uhlmann. Semiglobal boundary rigidity for Riemannian metrics, *Math. Ann.*, 325(2003), 767-793.
- [21] R. Michel. Sur la rigidité imposée par la longueur des géodésiques, *Invent. Math.*, 65(1981), 71-83.
- [22] R. G. Mukhometov. The reconstruction problem of a two-dimensional Riemannian metric, and integral geometry, *Soviet Math. Dokl*, 18(1), 27-37.

- [23] R. G. Mukhometov. On a problem of reconstructing Riemannian metrics, *Siberian Math. J.*, 22(3), 420-433.
- [24] R. G. Mukhometov and V. G. Romanov. On the problem of finding an isotropic Riemannian metric in an n -dimensional space (in Russian), *Dokl. Akad. Nauk SSSR* 243(1)(1978), 41-44.
- [25] J. P. Otal. Sur les longueurs des géodésiques d'une métrique à courbure négative dans le disque, *Comment. Math. Helv.* 65(1990), 334-347.
- [26] M. Porrati and R. Rabadan. Boundary rigidity and holography, *J. High. Energy Phys.*, 0401(2004), 034-057.
- [27] L. Pestov and G. Uhlmann. Two dimensional simple compact manifolds with boundary are boundary rigid, *Ann. of Math.*, 161(2)(2005), 1089-1106.
- [28] V. A. Sharafutdinov. *Integral Geometry of Tensor Fields*, VSP, Utrecht, the Netherlands, 1994.
- [29] V. A. Sharafutdinov. *Ray Transform on Riemannian Manifolds. Lecture Notes*, University of Oulu, 1999.
- [30] P. Stefanov and G. Uhlmann. Rigidity for metrics with the same lengths of geodesics, *Math. Res. Lett.* 8(1-2)(2001), 105-124.
- [31] P. Stefanov and G. Uhlmann. Stability estimates for the X-ray transform of tensor fields and boundary rigidity, *Duke Math. J.*, 123 (2004), 445-467.
- [32] P. Stefanov and G. Uhlmann. Boundary rigidity and stability for generic simple metrics, *Journal of AMS*, 18(2005), 975-1003.
- [33] P. Stefanov and G. Uhlmann. Integral geometry of tensor fields on a class of non-simple Riemannian manifolds, *Amer. J. of Math.*, 130 (2008), 239-268.
- [34] P. Stefanov and G. Uhlmann. Local lens rigidity with incomplete data for a class of non-simple Riemannian manifolds, *J. Diff. Geom.*, 82(2009), 383-409.
- [35] P. Stefanov and G. Uhlmann. The geodesic X-ray transform with fold caustics, *Anal. and PDE*, 5(2012), 219-260.
- [36] P. Stefanov, G. Uhlmann and A. Vasy. Boundary rigidity with partial data, *arXiv: 1306.2995*, 2013.
- [37] G. Uhlmann. The scattering relation and the Cauchy data, *Geometric methods in inverse problems and PDE control*, IMA publications, 137(2003), 263-288.
- [38] J. Vargo. A proof of lens rigidity in the category of analytic metrics, *Math. Res. Letters*, 16 (2009), 1057-1069.
- [39] H. Zhang. On the stability/sensitivity of recovering velocity fields from boundary measurements, (PhD dissertation), Michigan State University, 2013.